

Random perturbations of nonlinear parabolic systems

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August 2, 2011

Abstract

Several aspects of regularity theory for parabolic systems are investigated under the effect of random perturbations. The deterministic theory, when strict parabolicity is assumed, presents both classes of systems where all weak solutions are in fact more regular, and examples of systems with weak solutions which develop singularities in finite time. Our main result is the extension of a regularity result due to Kalita to the stochastic case. Concerning the examples with singular solutions (outside the setting of Kalita's regularity result), we do not know whether stochastic noise may prevent the emergence of singularities, as it happens for easier PDEs. We can only prove that, for a linear stochastic parabolic system with coefficients outside the previous regularity theory, the expected value of the solution is not singular.

MSC (2010): 60H15, 35B65, 35R60 (primary); 60H30 (secondary)

1 Introduction

Nonlinear parabolic systems of the form

$$\partial_t u = \operatorname{div} A(x, t, u, Du), \quad u|_{t=0} = u_0 \quad (1.1)$$

on a cylindrical domain $D \times (0, T)$, with $D \subset \mathbb{R}^n$ a bounded, regular domain, $u: D \times [0, T] \rightarrow \mathbb{R}^N$, $A: D \times [0, T] \times \mathbb{R}^N \times \mathbb{R}^{nN} \rightarrow \mathbb{R}^{nN}$, have been investigated by many authors, see for instance [14, 15, 13] and references therein. A key feature in the vectorial case $N > 1$ is that, under the strict parabolicity assumption

$$\sum_{i,j=1}^n \sum_{\alpha,\beta=1}^N \frac{\partial A_i^\alpha}{\partial z_j^\beta}(x, t, u, z) \xi_i^\alpha \xi_j^\beta \geq \lambda_0 |\xi|^2 \quad \text{for all } \xi \in \mathbb{R}^{nN}$$

and some differentiability assumptions with respect to the (x, u) -variable, there are classes of vector fields $A(x, t, u, z)$ such that all weak solutions to (1.1) are in fact more regular, and examples of systems such that there exist weak solutions with singularities; this dichotomy does not happen for single equations, the case $N = 1$, where regularity of weak solutions is always true, due to the (elliptic and parabolic) works based on the fundamental results of De Giorgi, Nash and Moser [4, 19, 18].

However, for second-order, parabolic systems under suitable additional assumptions on growth and regularity of the vector field $A(x, t, u, z)$, there are partial regularity results available, yielding Hölder regularity of the solution u (or of its spatial gradient Du) outside of a negligible set, the singular set of u (or of Du). Hence, for general systems which are nonlinear in the gradient variable, the best regularity to hope for is partial regularity of Du , with an estimate for the Hausdorff dimension of the singular set strictly below the dimension of $\mathbb{R}^n \times [0, T]$, see [7]. For better estimates on the Hausdorff dimension one needs to assume stronger assumptions (or also some a priori information on the regularity of the solution). Regularity of u on a larger set can be obtained, for instance, in low dimensions or with special structure assumptions (such as vector fields linear in the gradient variable), see e.g. [12, 2, 20]. Full regularity of u instead is only possible if even more restrictive structural assumptions are imposed. The easiest (and very classical one) of such

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examples are linear parabolic systems with constant coefficients. In the nonlinear case, in the famous case of the p -Laplacian system it is also possible to prove full regularity of Du , see [6]. Furthermore, if the system is still sufficiently close to the Laplacian system, then we still get full regularity of u , see [14, 13]. This regularity result will be of great importance for our paper.

Let us now go into some details, point out some of the structural prerequisites of the positive (full) regularity theory and confront it with the existing examples of systems admitting a singular weak solutions. For simplicity we focus here in the introduction on the case of quasilinear problems with a vector field of the form $A(x, t)z$, i. e. to weak solutions of

$$\partial_t u = \operatorname{div} (A(x, t) Du), \quad u|_{t=0} = u_0. \quad (1.2)$$

As mentioned above, without additional structural conditions on the coefficients full regularity of the solutions can no longer be expected in the vectorial case. It was observed by Koshelev and Kalita [14, 13] that if the coupling of the single equations is sufficiently weak, then discontinuities of the weak solution can globally be excluded:

Theorem 1.1 ([13]). *Let $u_0 \in W^{1,q}(D, \mathbb{R}^N)$ for some $q > n$ and consider coefficients $A(x, t)$ which are of class C^1 in x , measurable in t and which satisfy*

$$\lambda_0 |\xi|^2 \leq \langle A(x, t) \xi, \xi \rangle, \quad |A(x, t) \xi| \leq \lambda_1 |\xi|, \quad \text{and} \quad |D_x A(x, t)| \leq L \quad (1.3)$$

for all $\xi \in \mathbb{R}^{nN}$, $(x, t, z) \in D \times [0, T] \times \mathbb{R}^{nN}$ and some positive constants λ_0, λ_1, L . If in addition $\frac{\lambda_0}{\lambda_1} > 1 - \frac{2}{n}$ holds, then every weak solution $u: D \times [0, T] \rightarrow \mathbb{R}^N$ to the initial boundary value problem (1.2) is of class $C_{\text{loc}}^{0,\alpha}(D \times [0, T], \mathbb{R}^N)$ for some $\alpha > 0$.

It is important to mention that the original results apply to more general systems, possibly nonlinear in the gradient variable, provided that the vector field $A(x, t, u, z)$ is sufficiently close to a quasilinear situation with a small dispersion ratio. First, Koshelev proved the existence of a regular solution (which in the situation above is the unique one) by studying an approximation of the system such that its solutions are regular and converge in a suitable norm to a solution of the original system. Later Kalita achieved the regularity result for all solutions with a direct argument (and not as a consequence of a suitable approximating sequence), which essentially relies on Moser's iterative method [18].

Under weaker assumptions than in the previous theorem, such a global regularity result can no longer be expected. In fact, in a very similar setting the following example of a system was proposed by Stara and John [26] (actually, the example was constructed on the full space and the solution can be traced back in time $t \rightarrow -\infty$), which admits a solutions that starts from a regular – in particular Hölder continuous – initial data and develops a singularity in finite time in the interior of the parabolic cylinder.

Theorem 1.2 ([26]). *Let $n = N \geq 3$. There exist initial data $u_0 \in W^{1,2n}(B_1(0), \mathbb{R}^n)$ and measurable, symmetric coefficients $A \in L^\infty(B_1(0) \times [0, 1], \mathbb{R}^{n^2 \times n^2})$, which are elliptic and bounded in the sense of (1.3)_{1,2} for all $(x, t) \in B_1(0) \times [0, 1]$, such that at least one of the solutions to the initial problem (1.2) develops a discontinuity in the origin $x = 0$ as $t \nearrow 1$.*

The coefficients constructed in [26] have a dispersion ratio $\frac{\lambda_0}{\lambda_1} < 1 - \frac{2}{n}$ below the critical one investigated in [13] as well as a lower regularity with respect to the x -variable. For this reason we cannot exclude that the solution develops the discontinuity due to an interaction with the irregular coefficients. However, a comparison with the positive regularity results in the elliptic theory (the stationary case) suggests that the dispersion ratio $\frac{\lambda_0}{\lambda_1}$ plays an important role. Indeed, in the elliptic case no regularity in the x -variable is required, and a modification of De Giorgi's counterexample to full regularity shows sharpness of the (elliptic) condition on the dispersion ratio (see [15, Section 2.5]). Unfortunately, we didn't find further counterexamples in the literature which could give a similar complete picture in the case of parabolic systems.

The aim of this paper is to investigate parts of this theory under the effect of random perturbations. The final aim of our research project, in analogy with recent results proved for other equations, is to show that the regularity theory of parabolic systems is, under random perturbations, in some sense (of course only up to a certain degree of regularity) not worse than the deterministic one, and possibly better. As in

the deterministic case there is more than one approach to the analysis of these problems, so we restrict the attention here only to a few directions. More precisely, we want to show two results.

First, we study systems with Itô noise of the form

$$du = \operatorname{div} (A(x, t) Du) dt + H(Du) dB_t, \quad u|_{t=0} = u_0 \quad (1.4)$$

(for H Lipschitz) where $(B_t)_{t \geq 0}$ is a Brownian motion of suitable dimension, and we generalize to the stochastic case one of the regularity results of the deterministic theory, the work of Kalita [13] (which was displayed also above for the special case of a quasilinear system). The passage from deterministic to stochastic of Kalita's approach contains at least one non trivial detail which is rather new in the stochastic setting: the weak solution u we start with is not, a priori, the limit of a sequence of smooth solutions of approximating equations (for instance, due to the nonlinearity, classical mollifiers are difficult to implement; in another direction, in some cases solutions exist as limits of Galerkin or other types of approximations, but we here start with a weak solution which a priori has not been constructed in that way). Therefore, it is not clear how to perform differential calculus on u ; and Kalita's approach is heuristically based on an equation satisfied by second space derivatives of u . Therefore one has to use finite difference quotients in place of derivatives, a classical method in the deterministic setting, but not common in the stochastic case. This leads to a number of technical novelties. At the end, we reach a full extension of Kalita result to a quite general stochastic case, which includes in particular perturbations in form of additive or of multiplicative noise. In the quasi-linear model case we obtain – as a particular case of Theorem 5.1 – the following result (the precise definition of weak solution is given in Definition 2.2 below).

Theorem 1.3. *Let $u_0 \in W^{1,q}(D, \mathbb{R}^N)$ for some $q > n$. Consider coefficients $A(x, t)$ which are of class C^1 in x , measurable in t and which satisfy (1.3) with $\frac{\lambda_0}{\lambda_1} > 1 - \frac{2}{n}$, and assume that H is Lipschitz continuous with Lipschitz constant $L_H < L_H^*$ for some sufficiently small L_H^* depending only on n, λ_0 and λ_1 . Then there exists $\alpha > 0$ depending only on n, λ_0, λ_1 and q such that every weak solution $u: D \times [0, T] \times \Omega \rightarrow \mathbb{R}^N$ to the initial boundary value problem (1.4) is of class $C_{\text{loc}}^{0,\alpha}(D \times [0, T], \mathbb{R}^N)$ with probability 1.*

Second, we investigate for these systems the problem recently considered for other classes of PDEs, see [8], namely the possibility that it is precisely the noise which prevents the emergence of singularities. The aim of this project, that we have reached only partially until now, would be to prove that, under assumptions on the vector field $A(x, t, u, z)$ such that there exist weak solutions with singularities in the deterministic case, there are no more singularities if we add a suitable noise. We can only prove an intermediate but promising result. We consider linear stochastic systems with Stratonovich bilinear multiplicative noise of the form

$$du = \operatorname{div} (A(x, t) Du) dt + \sigma Du \circ dB_t, \quad u|_{t=0} = u_0 \quad (1.5)$$

with regular, bounded and elliptic measurable coefficients A . In this situation we obtain regularity for the mean value.

Proposition 1.4. *Given coefficients $A(x, t)$ which are of class C^1 in x , measurable in t , and which satisfy (1.3), there exists $\sigma_0 \geq 0$ depending only on λ_0, λ_1 such that for all $\sigma > \sigma_0$, all initial conditions $u_0 \in W^{1,q}(D, \mathbb{R}^N)$ with $q > n$, and all weak solutions u of equation (1.5) we have that the function $(x, t) \mapsto E[u(x, t)]$ is locally Hölder continuous on $D \times [0, T]$.*

This result is proved by applying the deterministic results. The key observation in the proof is that the average solves an equation with a better (that is greater) dispersion ratio. So far this class does not cover the counter-examples in [26] mentioned above (because here we need more regularity of A than in the counter-example), but however the theory presented here might be of its own interest. We cannot take $\sigma_0 = 0$ but we suspect that this is the critical value (namely that for all noise intensities the regularization takes place). The fact that this result holds independently of the initial condition u_0 – though sufficiently regular – and of the specific form of $A(x, t)$ in the functional class we consider (in particular the fact that no symmetry is embodied in our assumptions which could justify compensations due to the expected value), leads us to think that in fact $u(x, t)$ itself is Hölder continuous, but we do not have a proof of this conjecture.

Several problems remain open:

- (i) whether counter-examples exist also in the stochastic case under some conditions on A ;
- (ii) the regularity of $u(x, t)$ itself in the regularization-by-noise above and other related issues, such as the value of σ_0 and a generalization to the nonlinear case;
- (iii) the generalization of other deterministic approaches to regularity.

Concerning the existence of weak solutions, we could give a quite general result, but since it is related to the generalization of the approach of [14, 15] to regularity, we postpone it to a future work.

2 Setting and assumptions

Consider $n, n' \in \mathbb{N}$ with $n \geq 2$, $T > 0$, and $D \subset \mathbb{R}^n$ a (regular) bounded domain. Let (Ω, \mathcal{F}, P) be a complete probability space with a filtration $(\mathcal{F}_t)_{t \geq 0}$, and let $(B_t)_{t \geq 0}$ be a standard n' -dimensional Brownian motion. Let further $A: D \times [0, T] \times \mathbb{R}^N \times \mathbb{R}^{nN} \times \Omega \rightarrow \mathbb{R}^{nN}$ be a vector field satisfying the following properties:

- A is progressively measurable, i. e. for every $t \in [0, T]$ the restriction of A to $D \times [0, t] \times \mathbb{R}^N \times \mathbb{R}^{nN} \times \Omega \rightarrow \mathbb{R}^{nN}$ is $\mathcal{B}(D) \times \mathcal{B}([0, t]) \times \mathcal{B}(\mathbb{R}^N) \times \mathcal{B}(\mathbb{R}^{nN}) \times \mathcal{F}_t$ measurable;
- $A(x, t, u, z, \omega)$ (usually abbreviated by $A(x, t, u, z)$) is differentiable in x, u and z (with \mathcal{F}_t -adapted derivatives), and it satisfies for P -almost all $\omega \in \Omega$ the following assumptions concerning growth and ellipticity:

$$\left\{ \begin{array}{l} |A(x, t, u, z)| \leq L (|z| + |u|^{\frac{n+2}{n}} + f^{\frac{a}{2}}(x, t)) \\ |\xi - \kappa D_z A(x, t, u, z) \xi|^2 \leq (1 - \nu^2) |\xi|^2 \\ |D_u A(x, t, u, z)| \leq L (|z|^{\frac{2}{n+2}} + |u|^{\frac{2}{n}} + f(x, t)) \\ |D_x A(x, t, u, z)| \leq L (|z| + |u|^{\frac{n+2}{n}} + f^2(x, t)) \end{array} \right. \quad (2.1)$$

for all $(x, t) \in D \times [0, T]$, $u \in \mathbb{R}^N$ and $z, \xi \in \mathbb{R}^{nN}$, some constants $\kappa, \nu, L > 0$, and an \mathcal{F}_t -adapted process f which with probability one belongs to $L^a(D \times [0, T])$ for a fixed number $a > n + 2$.

Moreover, let $H: D \times [0, T] \times \mathbb{R}^{nN} \times \Omega \rightarrow \mathbb{R}^{nN}$ be progressively measurable, of class C^1 in x , Lipschitz with respect to the gradient variable of at most linear growth, uniformly in (x, t) , i.e.

$$\left\{ \begin{array}{l} |H(x, t, z, \omega) - H(x, t, \tilde{z}, \omega)| \leq L_H |z - \tilde{z}|, \\ |H(x, t, z, \omega)| \leq L (f_H(x, t, \omega) + |z|), \\ |D_x H(x, t, z, \omega)| \leq L (f_H^{\frac{a}{n+2}}(x, t, \omega) + |z|) \end{array} \right. \quad (2.2)$$

for a constant L_H , all $(x, t) \in D \times [0, T]$, $z, \tilde{z} \in \mathbb{R}^{nN}$, and almost every $\omega \in \Omega$. Here, f_H denotes a function in $L^a(D \times (0, T) \times \Omega)$.

Under these assumptions we consider a stochastic partial differential equation with noise of the form

$$du = \operatorname{div} A(x, t, u, Du) dt + H(x, t, Du) dB_t \quad \text{in } D_T := D \times (0, T), \quad (2.3)$$

where $u: D_T \times \Omega \rightarrow \mathbb{R}^N$ is a random function. The stochastic integral is here understood in the Itô sense. According to the growth condition on the vector field A , we note that for P -almost every $\omega \in \Omega$ and all $t \in [0, T]$ we have $\operatorname{div} A(x, t, v, Dv) \in W^{-1,2}(D, \mathbb{R}^N)$ – the dual space to $W_0^{1,2}(D, \mathbb{R}^N)$ –, provided that $v \in W^{1,2}(D, \mathbb{R}^N)$.

Remark 2.1. *We have chosen this level of generality of the noise for two reasons: to keep a simple PDE structure instead of an abstract operator formulation, and to cover two interesting examples: additive noise (with $H(x, z)$ independent of z) and bilinear multiplicative noise with first derivatives of u (with $H(x, z)$ linear in z). A priori there is no conceptual obstacle to consider H depending also on u and also to generalize to the case of a Brownian motion B in a Hilbert space U , with suitable assumptions on H , but for simplicity we restrict ourselves to the previous case.*

The function spaces that will be needed in the sequel are the Banach spaces

$$\begin{aligned} V^{m,p}(D_T, \mathbb{R}^N) &:= L^\infty(0, T; L^m(D, \mathbb{R}^N)) \cap L^p(0, T; W^{1,p}(D, \mathbb{R}^N)) \\ V_0^{m,p}(D_T, \mathbb{R}^N) &:= L^\infty(0, T; L^m(D, \mathbb{R}^N)) \cap L^p(0, T; W_0^{1,p}(D, \mathbb{R}^N)), \end{aligned}$$

with $m, p \geq 1$, and they are equipped with the norm

$$\|u\|_{V^{m,p}(D_T, \mathbb{R}^N)} := \operatorname{ess\,sup}_{t \in (0, T)} \|u(t)\|_{L^m(D, \mathbb{R}^N)} + \|Du\|_{L^p(D_T, \mathbb{R}^N)}.$$

When $m = p$ we shall use the abbreviations $V_{(0)}^p(D_T, \mathbb{R}^N) = V_{(0)}^{p,p}(D_T, \mathbb{R}^N)$. We remind that the spaces $V^{m,p}(D_T, \mathbb{R}^N)$ are embedding in the Lebesgue space $L^q(D_T, \mathbb{R}^N)$ with $q = p^{\frac{n+m}{n}} > p$ (see [6, Propositions I.3.1, I.3.2]). We will need only the result concerning the cases $m = 2$, $p \geq 2$ or $m = p \geq 2$. In the latter case, the embedding reads as follows (see [6, Propositions I.3.3, I.3.4]): let $v \in V_0^p(D_T, \mathbb{R}^N)$, $p < n$. Then there exists a constant c depending only on n and p such that

$$\|v\|_{L^q(D_T, \mathbb{R}^N)} \leq c \|v\|_{V^p(D_T, \mathbb{R}^N)} \quad (2.4)$$

(and an analogous result holds without any restriction on the boundary values of v on $\partial D \times (0, T)$ if ∂D is assumed to be sufficiently regular).

We are now going to study the properties of weak (or variational) solutions to the system (2.3), which is to be understood in the following sense.

Definition 2.2. *An \mathcal{F}_t -progressively measurable process u on $[0, T] \times \Omega$ is called a weak solution to the system (2.3) with initial values $u_0 \in L^2(D, \mathbb{R}^N)$ if P -a. e. path satisfies $u(\cdot, \omega) \in V^2(D_T, \mathbb{R}^N)$ and if for all $t \in [0, T]$, we have P -a. s. the identity*

$$\langle u(t) - u_0, \varphi \rangle_{L^2(D)} = \int_0^t \langle \operatorname{div} A(\cdot, s, u, Du), \varphi \rangle_{W^{-1,2}(D); W_0^{1,2}(D)} ds + \int_0^t \langle \varphi, H(\cdot, s, Du) dB_s \rangle_{L^2(D)}$$

for all $\varphi \in W_0^{1,2}(D, \mathbb{R}^N)$.

When a solution is progressively measurable with respect to the (completed) filtration associated to the Brownian motion, it is usually called a “strong” solution in the probabilistic sense, see [23, Section IX.1]. We do not require this condition, so our result will also apply to the so called “weak” solutions in the probabilistic sense (those for which there is a filtration $(\mathcal{F}_t)_{t \geq 0}$ such that u is \mathcal{F}_t -progressively measurable and B is an \mathcal{F}_t -Brownian motion). We further note that according to the definition above, a solution is defined as an equivalence class in the sense of versions (a process Y is a version or modification of a process X if for each time t we have P -a. s. $X_t = Y_t$). Hence, regularity of a weak solution is always to be understood as finding a regular representative in the corresponding equivalence class.

Moreover, we comment on the way in which the initial values are attained. Under mild assumptions on the growth of A and H with respect to the gradient variable one actually deduces from the equation itself that u belongs to $C^0(0, T; L^2(D', \mathbb{R}^N))$ P -a. s. for every $D' \Subset D$, compare formula (4.1) and the beginning of Step 3 on p. 15. Under further assumptions on the trace of u on $\partial D \times [0, T]$ this extends to continuity of the full L^2 -norm, with $D' = D$. In this sense the term “initial value” in the definition of a weak solution as a function in the space V^2 is justified.

3 Preliminaries

In this section we recall some well-known facts and provide some technical tools. For convenience of the reader we state two suitable versions of Itô’s formula. Furthermore, in analogy with the deterministic theory, we discuss a sufficient condition for the “existence of weak derivatives with probability one”, and we further give a criterion which guarantees pathwise Hölder continuity.

3.1 Itô formula

We first recall two versions of Itô's formula, the first one the standard version for N -dimensional processes and the second one for processes with values in Hilbert spaces. Consider (Ω, \mathcal{F}, P) a complete probability space and let

$$dX(t) = a(t) dt + b(t) dB_t \quad (3.1)$$

be an N -dimensional Itô process which satisfies: a, b are \mathcal{F}_t -adapted (i.e., the maps $\omega \mapsto a(t, \omega), b(t, \omega)$ are \mathcal{F}_t measurable), $(t, \omega) \mapsto b(t, \omega)$ is $\mathcal{B}([0, T]) \times \mathcal{F}$ -measurable and

$$P\left(\int_0^T [|a(s, \omega)| + |b(s, \omega)|^2] ds < \infty\right) = 1.$$

Then the following general Itô formula holds (see e.g. [21, Theorem 4.2.1]).

Theorem 3.1 (Itô's formula I). *Let $g(t, z) = (g_1(t, z), \dots, g_p(t, z))$ be a map from $[0, T] \times \mathbb{R}^N$ to \mathbb{R}^p of class C^1 in t and of class C^2 in z . Then the process $Y(t, \omega) := g(t, X(t))$ with $X(t)$ defined in (3.1) is again an Itô process whose components are given by*

$$dY_k(t) = \frac{\partial g_k}{\partial t}(t, X) dt + \sum_{i=1}^N \frac{\partial g_k}{\partial y_i}(t, X) dX_i + \frac{1}{2} \sum_{i,j=1}^N \frac{\partial^2 g_k}{\partial y_i \partial y_j}(t, X) dX_i dX_j$$

for all $k \in \{1, \dots, p\}$, and with $dB_i dB_j = \delta_{ij} dt$ and $dB_i dt = 0 = dt dB_i$ for all $i, j \in \{1, \dots, N\}$.

In the sequel, we will also employ the following version of the Itô formula in Hilbert spaces that can be found in [16, Theorem 3.1] or [24, Chapter 4.2, Theorem 2].

Theorem 3.2 (Itô's formula II). *Let $V \subset H \subset V'$ be a Gelfand triple, with H a separable Hilbert space. Assume that we have for $\mathcal{L}^1 \times P$ almost all (t, ω)*

$$\langle x(t), \varphi \rangle_H = \langle x(0), \varphi \rangle_H + \int_0^t \langle y(s), \varphi \rangle_{V', V} ds + \langle M_t, \varphi \rangle_H \quad (3.2)$$

for every $\varphi \in V$ where $x(t, \omega), y(t, \omega)$ are taking values in V and V' , respectively, and are progressively measurable with

$$P\left(\int_0^T [\|x(s, \omega)\|_V^2 + \|y(s, \omega)\|_{V'}^2] ds < \infty\right) = 1,$$

and where M_t is a continuous local martingale with values in H . Then there exists a set $\tilde{\Omega} \subset \Omega$ with $P(\tilde{\Omega}) = 1$ and a map $\tilde{x}(t, \omega)$ with values in H such that:

(i) $\tilde{x}(t)$ is \mathcal{F}_t -adapted, continuous in $t \in [0, T]$ for every $\omega \in \tilde{\Omega}$, and $x(t) = \tilde{x}(t)$ P -almost surely;

(ii) for every $\omega \in \tilde{\Omega}$ and every $\varphi \in V$ there holds

$$\langle \tilde{x}(t), \varphi \rangle_H = \langle x(0), \varphi \rangle_H + \int_0^t \langle y(s), \varphi \rangle_{V', V} ds + \langle M_t, \varphi \rangle_H;$$

(iii) for every $\omega \in \tilde{\Omega}$ there holds the inequality

$$\|\tilde{x}(t)\|_H^2 = \|x(0)\|_H^2 + 2 \int_0^t \langle y(s), x(s) \rangle_{V', V} ds + 2 \int_0^t \langle dM_s, \tilde{x}(s) \rangle_H + [M]_t.$$

3.2 Weak derivatives

For a vector-valued function $f: \mathbb{R}^n \supset D \rightarrow \mathbb{R}^N$, $k \in \{1, \dots, n\}$ and a real number $h \in \mathbb{R}$ we denote by $\Delta_{k,h}f(x) := h^{-1}(f(x + he_k) - f(x))$ the finite different quotient in direction e_k and stepsize h (this makes sense as long as $x, x + he_k \in D$). Let $p > 1$, $f \in L^p(D)$, $k \in \{1, \dots, n\}$ and let $D_k f$ be the derivative of f in the direction k in the sense of distributions. Just for comparison, let us recall the following lemma (not used below).

Lemma 3.3. *If there is $h_n \rightarrow 0$ and $g_k \in L^p(D)$ such that*

$$\lim_{n \rightarrow \infty} \int_D (\Delta_{k,h_n} f(x) - g_k(x)) \varphi(x) dx = 0$$

for every $\varphi \in C_0^\infty(D)$, then $D_k f$ is in $L^p(D)$ and is equal to g_k .

As an immediate consequence of this lemma and of the compactness of the L^p -spaces with $p > 1$ with respect to weak (or weak-*) convergence, we obtain a simple criterion for the existence of the weak derivative $D_k f$ in L^p , namely it is sufficient that $\|\Delta_{k,h} f\|_{L^p(D')}$ is bounded by some constant, for all h and every $D' \Subset D$ such that $|h| < \text{dist}(D', \partial D)$.

Now this well-known principle shall be carried over to a probabilistic setting. Let (Ω, F, P) be a complete probability space and consider a function f in the Banach space $L^p(D \times \Omega)$. A function $g_k \in L^p(D \times \Omega)$ is said to be the weak derivative of f in the k -direction if

$$P\left(\int_D f D_k \varphi dx = - \int_D g_k \varphi dx\right) = 1$$

for every $\varphi \in C_0^\infty(D)$ (taking a countable sequence and using a density argument, the property “for every $\varphi \in C_0^\infty(D)$ ” can be written inside the probability). We then write $D_k f = g_k$. The previous lemma has a generalization to functions in $L^p(D \times \Omega)$.

Lemma 3.4. *If there is $h_n \rightarrow 0$ and $g_k \in L^p(D \times \Omega)$ such that*

$$\lim_{n \rightarrow \infty} \int \int_{D \times \Omega} (\Delta_{k,h_n} f(x, \omega) - g_k(x, \omega)) \varphi(x) X(\omega) dx dP(\omega) = 0$$

for every $\varphi \in C_0^\infty(D)$ and every bounded measurable $X: \Omega \rightarrow \mathbb{R}$, then $D_k f$ is in $L^p(D \times \Omega)$ and is equal to g_k .

Proof. Since X and φ are bounded, we may apply (first Fubini and then) Lebesgue’s dominated convergence theorem, and we get

$$\begin{aligned} -E\left[X \int_D (f D_k \varphi + g_k \varphi) dx\right] &= - \int \int_{D \times \Omega} X (f D_k \varphi + g_k \varphi) dx dP \\ &= \lim_{n \rightarrow \infty} \int \int_{D \times \Omega} X (-f(x) \Delta_{k,-h_n} \varphi(x) - g_k(x) \varphi(x)) dx dP. \end{aligned}$$

When $h_n < \text{dist}(\text{spt } \varphi, \partial D)$, this is equal to (we apply Fubini twice and a change of variables)

$$\lim_{n \rightarrow \infty} \int \int_{D \times \Omega} X (\Delta_{k,h_n} f(x, \omega) \varphi(x) - g_k(x, \omega) \varphi(x)) dx dP.$$

This limit is zero by assumption, hence

$$E\left[X \int_D (f D_k \varphi + g_k \varphi) dx\right] = 0.$$

The arbitrariness of X implies $\int_D (f D_k \varphi + g_k \varphi) dx = 0$, as a random variable on Ω . The proof is complete. \square

Corollary 3.5. *If there is a constant $C > 0$ such that*

$$E \left[\int_{D'} |\Delta_{k,h} f(x)|^p dx \right] \leq C$$

for all h and all $D' \Subset D$ such that $|h| < \text{dist}(D', \partial D)$, then $D_k f$ is in $L^p(D \times \Omega)$.

Proof. The family $g_{k,h}(x, \omega) := \Delta_{k,h} f(x, \omega)$ is equibounded in $L^p(D' \times \Omega)$, hence there is a sequence $h_n \rightarrow 0$ such that g_{k,h_n} converges weakly in $L^p(D \times \Omega)$ to some function $g_k \in L^p(D \times \Omega)$. The product $\varphi(x)X(\omega)$ is in $L^{p'}(D \times \Omega)$ (p' conjugate to p) for every $\varphi \in C_0^\infty(D)$ and every bounded measurable $X : \Omega \rightarrow \mathbb{R}$. Hence, we may apply the lemma and obtain the assertion. \square

First, for our later application, we replace D by $D \times [0, T]$ and we allow different integrability exponents with respect to the variables in $[0, T]$ and D , respectively. Let $f : \Omega \rightarrow L^q(0, T; L^p(D))$ be a measurable function with $p \in (1, \infty)$ and $q > 1$. We say that a function $g_k : \Omega \rightarrow L^q(0, T; L^p(D))$ is weak derivative of f in the k -direction with probability one if for a.e. $(t, \omega) \in [0, T] \times \Omega$ we have

$$\int_D f D_k \varphi dx = - \int_D g_k \varphi dx$$

for every $\varphi \in C_0^\infty(D)$, and we then write $D_k f = g_k$. Furthermore, let us generalize to a scheme where we relax the integrability in Ω .

Theorem 3.6. *Let $Y : [0, T] \times \Omega \rightarrow (0, 1]$ be a positive random variable, with $P(\inf_{t \in [0, T]} Y > 0) = 1$. If there is a constant $C > 0$ such that*

$$E \left[\|Y(t) \Delta_{k,h_n} f(x, t)\|_{L^q(0, T; L^p(D'))}^p \right] \leq C$$

for all h and $D' \Subset D$ satisfying $|h| < \text{dist}(D', \partial D)$, then $D_k f \in L^q(0, T; L^p(D))$ with probability one and there hold

$$\begin{aligned} Y \Delta_{k,h} f &\rightarrow Y D_k f \quad \text{weakly in } L^p(\Omega; L^q(0, T; L^p(D))), \\ E \left[\|Y D_k f\|_{L^q(0, T; L^p(D))}^p \right] &\leq C \end{aligned}$$

with the same constant C .

Proof. The family $Z_{k,h}(x, t, \omega) := Y(t) \Delta_{k,h_n} f(x, t, \omega)$ is equibounded in $L^p(\Omega; L^q(0, T; L^p(D')))$, hence there is a sequence $h_n \rightarrow 0$ such that Z_{k,h_n} converges weakly in $L^p(\Omega; L^q(0, T; L^p(D))$ (or weakly-* if $q = \infty$) to some function $Z_k \in L^p(\Omega; L^q(0, T; L^p(D))$). This implies (again with ψ, X bounded, measurable and φ smooth, compactly supported)

$$\lim_{n \rightarrow \infty} \int \int \int_{D \times [0, T] \times \Omega} (Y(t) \Delta_{k,h_n} f(x, t) - Z_k(x, t)) \varphi(x) \psi(t) X dx dt dP = 0.$$

Hence, by Fubini and change of variables as above, we find

$$\lim_{n \rightarrow \infty} \int \int \int_{D \times [0, T] \times \Omega} (Y(t) f(x, t) \Delta_{k,-h_n} \varphi(x) + Z_k(x, t) \varphi(x)) \psi(t) X dx dt dP = 0,$$

which in turn implies by Lebesgue's theorem

$$\int \int \int_{D \times [0, T] \times \Omega} (Y(t) f(x, t) D_k \varphi(x) + Z_k(x, t) \varphi(x)) \psi(t) X dx dt dP = 0.$$

Arbitrariness of X and ψ thus yields

$$\int_D (Y(t) f(x, t) D_k \varphi(x) + Z_k(x, t) \varphi(x)) dx = 0$$

for a. e. $(t, \omega) \in [0, T] \times \Omega$. Therefore, we have

$$\int_D (f(x, t) D_k \varphi(x) + g_k(x, t) \varphi(x)) dx = 0 \quad (3.3)$$

for a. e. $(t, \omega) \in [0, T] \times \Omega$, where $g_k = Y^{-1} Z_k$. Since Z_k belongs to $L^p(\Omega; L^q(0, T; L^p(D)))$, it is $L^q(0, T; L^p(D))$ for P -a. e. $\omega \in \Omega$. Hence, by assumption on Y , we also have $g_k \in L^q(0, T; L^p(D))$ for P -a. e. $\omega \in \Omega$. The only difference with the definition of g_k being the “weak derivative of f in the k -direction with probability one” is that the negligible set of $(t, \omega) \in [0, T] \times \Omega$ where (3.3) may fail depends on $\varphi \in C_0^\infty(D)$, until now. But $W^{1,p'}(D)$ (p' is conjugate to p) is separable and $C_0^\infty(D)$ is dense in it. Hence, there is a countable family $\{\varphi_n\} \subset C_0^\infty(D)$ which is dense in $W^{1,p'}(D)$. If we call N the countable union of all negligible sets of $(t, \omega) \in [0, T] \times \Omega$ where (3.3) may fail for $\{\varphi_n\}$, N is negligible, and on the complementary we have (3.3) for every φ_n , hence by density for all $\varphi \in W^{1,p'}(D)$ and then for all $\varphi \in C_0^\infty(D)$. Having identified g_k as the weak derivative of f in the k -direction we take advantage of the lower semi-continuity of the norm with respect to weak (or weak-*) convergence and thus find

$$E[\|Y D_k f\|_{L^q(0,T;L^p(D))}^p] = E[\|Z_k\|_{L^q(0,T;L^p(D))}] \leq C.$$

The proof is complete. \square

Remark 3.7. This result will be applied later in the cases $p = q$ where the assumption then reads as

$$E\left[\int_0^T \int_D |Y(t) \Delta_{k,h} f_1(x, t)|^p dx dt\right] \leq C$$

and where we have $L^p(\Omega; L^q(0, T; L^p(D))) = L^p(D \times [0, T] \times \Omega)$, or in the case $q = \infty$ where we then require

$$E\left[\sup_{t \in (0, T)} \int_D |Y(t) \Delta_{k,h} f_2(x, t)|^p dx\right] \leq C.$$

From the theorem we then conclude that $D_k f_1 \in L^p(D \times [0, T])$ and $D_k f_2 \in L^\infty(0, T; L^p(D))$ with probability one, respectively. In particular, if we take a function $f \in W^{1,p}(D)$ and if the previous assumptions are satisfied for $f_1 = Df$ and $f_2 = f$, then the conclusions are equivalent to $D_k f \in V^p(D_T)$.

3.3 A criterion for pathwise Hölder continuity

We next discuss a criterion which guarantees Hölder continuity of (a suitable representative of) a given functions $u: D \times [0, T] \rightarrow \mathbb{R}^N$. For example, Sobolev's embedding theorem provides a criterion easy to apply if u is in a suitable Sobolev space $W^{1,q}(D \times [0, T], \mathbb{R}^N)$ – but which in general is not satisfied for the solutions considered in our paper since derivatives in time need not exist. Instead, we now prove that it is sufficient that only the spatial derivatives belong to a suitable Lebesgue space, provided that a weak form of continuity in time (i. e. of the $L^2(D)$ -norm) is available.

Lemma 3.8. If a function $u: D \times [0, T] \rightarrow \mathbb{R}^N$ has the properties

$$Du \in L^\infty(0, T; L^{n+\alpha}(D, \mathbb{R}^N)), \quad u \in C^\beta(0, T; L^2(D, \mathbb{R}^N))$$

for some $\alpha, \beta > 0$, $D \subset \mathbb{R}^n$ a bounded, regular domain, then

$$u \in C^\gamma(D \times [0, T], \mathbb{R}^N)$$

for some $\gamma > 0$, depending only on α, β and n .

Proof. First, we deduce spatial Hölder continuity for every time slice. From the assumption $Du \in L^\infty(0, T; L^{n+\alpha}(D))$ we deduce $u \in L^\infty(0, T; C^\delta(D))$ for some $\delta > 0$, depending only on α and n , by Sobolev's embedding theorem. Namely, there exists $C_1 > 0$ such that

$$|u(x, t) - u(y, t)| \leq C_1 |x - y|^\delta \quad (3.4)$$

for all $t \in [0, T]$, $x, y \in D$.

Our next aim is Hölder continuity in time, at a fixed point. From the inequality

$$\|u(\cdot, t) - u(\cdot, s)\|_{L^2(D)} \leq C_2 |t - s|^\beta$$

for $s, t \in [0, T]$, we infer for every set $B \subset D$

$$\inf_{x \in B} |u(x, t) - u(x, s)| \leq \frac{1}{|B|} \int_B |u(x, t) - u(x, s)| dx \leq \frac{1}{|B|^{1/2}} \|u(\cdot, t) - u(\cdot, s)\|_{L^2(D)} \leq \frac{C_2 |t - s|^\beta}{|B|^{1/2}}.$$

Let $x_0 \in D$ be given. In order to prove Hölder continuity in time at x_0 , we estimate

$$\begin{aligned} |u(x_0, t) - u(x_0, s)| &\leq |u(x_0, t) - u(x, t)| + |u(x, t) - u(x, s)| + |u(x, s) - u(x_0, s)| \\ &\leq 2C_1 |x - x_0|^\delta + |u(x, t) - u(x, s)| \end{aligned}$$

for every $x \in D$. Hence, if we take x in a ball $B(x_0, \rho)$, we have

$$\begin{aligned} |u(x_0, t) - u(x_0, s)| &\leq 2C_1 \rho^\delta + \inf_{x \in B(x_0, \rho)} |u(x, t) - u(x, s)| \\ &\leq 2C_1 \rho^\delta + C_3 \frac{C_2 |t - s|^\beta}{\rho^{n/2}} \end{aligned}$$

where C_3 is such that $|B(x_0, \rho)| = \rho^n / C_3^2$. Let us now choose $\rho = |t - s|^\varepsilon$ for some $\varepsilon > 0$:

$$|u(x_0, t) - u(x_0, s)| \leq 2C_1 |t - s|^{\varepsilon\delta} + C_3 C_2 |t - s|^{\beta - \varepsilon n/2}.$$

If we choose for instance $\varepsilon = \beta/n$, we get

$$|u(x_0, t) - u(x_0, s)| \leq C_4 |t - s|^\eta \tag{3.5}$$

for some $\eta, C_4 > 0$, independently of $x_0 \in D$, $t, s \in [0, T]$. The constant η depends only on β, δ and n .

From (3.4) and (3.5) it is now straightforward to deduce the claim of the lemma. \square

With the previous lemma at hand, we now give a criterion in the probabilistic setting, with (Ω, F, P) a complete probability space, which is adapted to weak solutions.

Proposition 3.9. *Let $u : D \times [0, T] \times \Omega \rightarrow \mathbb{R}^{nN}$ have the properties*

$$P(Du \in L^\infty(0, T; L^{n+\varepsilon}(D, \mathbb{R}^{nN}))) = 1 \tag{3.6}$$

$$u(x, t) = u_0(x) + \int_0^t a(x, s) ds + \int_0^t b(x, s) dB_s$$

for some $\varepsilon > 0$, $u_0 \in L^2(D)$, and with progressively measurable fields a, b such that

$$P\left(\int_0^T \int_D |a(x, s)|^2 dx ds + \int_0^T \left(\int_D |b(x, s)|^2 dx\right)^{\frac{2+\varepsilon}{2}} ds < \infty\right) = 1.$$

Then

$$P(u \in C^\gamma(D \times [0, T])) = 1$$

for some $\gamma > 0$ depending only on ε .

Proof. Step 1. If we prove that, for some $\beta > 0$,

$$P(u \in C^\beta(0, T; L^2(D))) = 1,$$

then we get the claim of the proposition after the pathwise application of the previous Lemma 3.8 (using in particular the stated independence of the Hölder exponent). To this end we observe that the function u is the sum of two terms:

$$u_1(x, t) = u_0(x) + \int_0^t a(x, s) ds, \quad u_2(x, t) = \int_0^t b(x, s) dB_s$$

The term u_1 is, with probability one, of class $W^{1,2}(0, T; L^2(D))$, hence it is of class $C^{1/2}(0, T; L^2(D))$:

$$\|u_1(t) - u_1(s)\|_{L^2(D)} = \left\| \int_s^t a(\cdot, r) dr \right\|_{L^2(D)} \leq |t - s|^{1/2} \left(\int_0^T \int_D |a(x, r)|^2 dx dr \right)^{1/2}.$$

So it only remains to prove that, for some $\beta > 0$,

$$P(u_2 \in C^\beta(0, T; L^2(D))) = 1.$$

Step 2. For $R > 0$, let

$$\tau_R = \inf \{t \in (0, T] : \int_0^t \|b(\cdot, s)\|_{L^2(D)}^{2+\varepsilon} ds > R\}$$

if the set is non empty, otherwise $\tau_R = T$. Let $\Omega_R \subset \Omega$ be the set where $\tau_R = T$. The family $\{\Omega_R\}_{R>0}$ is increasing, with

$$P\left(\bigcup_{R>0} \Omega_R\right) = 1$$

because by assumption we have $P(\int_0^T \|b(\cdot, s)\|_{L^2(D)}^{2+\varepsilon} ds < \infty) = 1$. We now set

$$b_R(x, s) = b(x, s)1_{s \leq \tau_R} \quad \text{and} \quad u_{2,R}(t) = \int_0^t b_R(x, s) dB_s = \int_0^{t \wedge \tau_R} b(x, s) dB_s.$$

We then have

$$\int_0^T \|b_R(\cdot, s, \omega)\|_{L^2(D)}^{2+\varepsilon} ds \leq R$$

uniformly in ω . Hence, for every $p \geq 1$, we find

$$\begin{aligned} E\left[\|u_{2,R}(t) - u_{2,R}(s)\|_{L^2(D)}^p\right] &= E\left[\left\|\int_s^t b_R(\cdot, r) dB_r\right\|_{L^2(D)}^p\right] \\ &\leq C_p E\left[\left(\int_s^t \|b_R(\cdot, r)\|_{L^2(D)}^2 dr\right)^{\frac{p}{2}}\right] \\ &\leq C_p |t - s|^{\frac{p\varepsilon}{2(2+\varepsilon)}} E\left[\left(\int_s^t \|b_R(\cdot, r)\|_{L^2(D)}^{2+\varepsilon} dr\right)^{\frac{p}{2+\varepsilon}}\right] \leq C_p R^{\frac{p}{2+\varepsilon}} |t - s|^{\frac{p\varepsilon}{2(2+\varepsilon)}}. \end{aligned}$$

This implies, for $p = p(\varepsilon)$ sufficiently large, by Kolmogorov's regularity theorem for processes taking values in $L^2(D)$ (see [3, Theorem 3.3] for a version in Banach spaces), that $u_{2,R}$ has a Hölder continuous version in $L^2(D)$

$$\|u_{2,R}(\cdot, t, \omega) - u_{2,R}(\cdot, s, \omega)\|_{L^2(D)} \leq C_{\beta,R}(\omega) |t - s|^\beta$$

with β any Hölder exponent with $\beta < \frac{\varepsilon}{2(2+\varepsilon)}$. For $\omega \in \Omega_R$ we thus have (recalling the definition of $u_{2,R}$)

$$\|u_2(\cdot, t, \omega) - u_2(\cdot, s, \omega)\|_{L^2(D)} \leq C_{\beta,R}(\omega) |t - s|^\beta.$$

Since $\bigcup_{R>0} \Omega_R$ is of full P -measure, we obtain $u_2 \in C^\beta(0, T; L^2(D))$ for P -a.e. $\omega \in \Omega$. Now the previous Lemma 3.8 can be applied, and the proof is complete. \square

3.4 A technical lemma

In Kalita's paper a crucial point is to show higher regularity (such as higher integrability and differentiability) not only for the solution, but also for powers of the solution (resp. its gradient). For this purpose the following technical lemma was essential.

Lemma 3.10 ([13]). *Let $u: \mathbb{R}^n \rightarrow \mathbb{R}^N$ be a function which is a.e. differentiable. Set $v = u|u|^s$ with $s \in (-1, \infty)$. Then, for $\mu(s) := 1 - (\frac{s}{2+s})^2$, we have a.e.*

$$Du \cdot Dv \geq \mu^{\frac{1}{2}}(s) |Du| |Dv|.$$

We need the following modification of this result, which on the one hand allows to test the system with powers (truncated for large values) and which on the other hand satisfies an estimate corresponding to the one from Lemma 3.10.

Lemma 3.11. *For every $K > 0$ and every $q \geq 1$ there exists a C^2 -function $T_{q,K}: \mathbb{R}^+ \rightarrow \mathbb{R}^+$ such that*

(i) $T_{q,K}$ is strictly increasing and convex on \mathbb{R}^+ , and it satisfies $T_{q,K}(t) = t^{2q}$ for all $t \leq K$;

(ii) for all $t \in \mathbb{R}^+$ and a constant $c(q)$ the growth with respect to t is estimated by

$$T_{q,K}(t) + T'_{q,K}(t)t + T''_{q,K}(t)t^2 \leq c(q) \min \{K^{2q-2}t^2, t^{2q}\};$$

moreover, the inequalities $T''_{q,K}(t)t - T'_{q,K}(t) \leq 2(q-1)T'_{q,K}(t)$ as well as $T''_{q,K}(t)t^2 \leq c(q)T'_{q,K}(t)t \leq c(q)T_{q,K}(t)$ hold true on \mathbb{R}^+ ;

(iii) If $u: \mathbb{R}^n \rightarrow \mathbb{R}^N$ is a function which is a.e. differentiable and $\mu(q) := 1 - (\frac{q-1}{q})^2$, then for the function $v = T'_{q,K}(|u|)|u|^{-1}u$ the following inequality is satisfied a.e.:

$$Du \cdot Dv \geq \sqrt{\mu(q)} |Du| |Dv| \geq \sqrt{\mu(q)} T'_{q,K}(|u|) |u|^{-1} |Du|^2.$$

Proof. We first assume $K = 1$. We set

$$T_{q,1}(t) = \begin{cases} t^{2q} & \text{if } t \leq 1 \\ at^2 + bt + c & \text{if } t > 1 \end{cases}$$

for some coefficients $a, b, c \in \mathbb{R}$ to be determined as follows. The C^2 -regularity condition implies that the following linear system has to be satisfied:

$$\begin{pmatrix} 1 & 1 & 1 \\ 2 & 1 & 0 \\ 2 & 0 & 0 \end{pmatrix} \begin{pmatrix} a \\ b \\ c \end{pmatrix} = \begin{pmatrix} 1 \\ 2q \\ 2q(2q-1) \end{pmatrix} \Rightarrow \begin{pmatrix} a \\ b \\ c \end{pmatrix} = \begin{pmatrix} q(2q-1) \\ -4q(q-1) \\ 1-3q+2q^2 \end{pmatrix}.$$

We now calculate some crucial quantities. We first observe that

$$T'_{q,1}(t) = \begin{cases} 2qt^{2q-1} & \text{if } t \leq 1 \\ 2q(2q-1)t - 4q(q-1) & \text{if } t > 1 \end{cases}$$

is strictly increasing and positive on \mathbb{R}^+ . Thus, we immediately obtain assertion (i) of the lemma. Furthermore, we have

$$T''_{q,1}(t)t - T'_{q,1}(t) = \begin{cases} 4q(q-1)t^{2q-1} & \text{if } t \leq 1 \\ 4q(q-1) & \text{if } t > 1, \end{cases}$$

which is again positive on \mathbb{R}^+ . Moreover, for all $t \in \mathbb{R}^+$ we obtain

$$T''_{q,1}(t)t - T'_{q,1}(t) \leq 2(q-1)T'_{q,1}(t), \quad (3.7)$$

which in particular yields the inequality $T''_{q,1}(t)t^2 \leq c(q)T'_{q,1}(t)t$ of assertion (ii). The last inequality $T'_{q,1}(t)t \leq c(q)T_{q,1}(t)$ is also checked easily. For the function v given in (iii) we next compute

$$\begin{aligned} D_i v^\alpha &= T'_{q,1}(|u|) \frac{D_i u^\alpha}{|u|} + (T''_{q,1}(|u|)|u| - T'_{q,1}(|u|)) \frac{D_i u \cdot u u^\alpha}{|u|^3}, \\ Du \cdot Dv &= T'_{q,1}(|u|) \frac{|Du|^2}{|u|} + (T''_{q,1}(|u|)|u| - T'_{q,1}(|u|)) \frac{|Du \cdot u|^2}{|u|^3}. \end{aligned}$$

In particular, this shows $Du \cdot Dv \geq 0$, using again the positivity of $T''_{q,1}(t)t - T'_{q,1}(t)$ and of $T'_{q,1}(t)$ on \mathbb{R}^+ . Furthermore, we obtain

$$\begin{aligned} |Dv|^2 &= T'_{q,1}(|u|)^2 \frac{|Du|^2}{|u|^2} + (T''_{q,1}(|u|)|u| - T'_{q,1}(|u|))^2 \frac{|Du \cdot u|^2}{|u|^4} \\ &\quad + 2 T'_{q,1}(|u|) (T''_{q,1}(|u|)|u| - T'_{q,1}(|u|)) \frac{|Du \cdot u|^2}{|u|^4} \geq T'_{q,1}(|u|)^2 \frac{|Du|^2}{|u|^2}, \end{aligned}$$

which yields the second inequality in (iii). Now, keeping in mind the definition of $\mu(\cdot)$, we find via the previous estimate (3.7)

$$\begin{aligned} |Du \cdot Dv|^2 - \mu(q) |Du|^2 |Dv|^2 &= \left(\frac{q-1}{q}\right)^2 T'_{q,1}(|u|)^2 \frac{|Du|^4}{|u|^2} + (T''_{q,1}(|u|)|u| - T'_{q,1}(|u|))^2 \frac{|Du \cdot u|^4}{|u|^6} \\ &\quad - (T''_{q,1}(|u|)|u| - T'_{q,1}(|u|))^2 \frac{2q-1}{q^2} \frac{|Du \cdot u|^2 |Du|^2}{|u|^4} \\ &\quad + 2 \left(\frac{q-1}{q}\right)^2 T'_{q,1}(|u|) (T''_{q,1}(|u|)|u| - T'_{q,1}(|u|)) \frac{|Du \cdot u|^2 |Du|^2}{|u|^4} \\ &\geq \left(\frac{q-1}{q}\right)^2 T'_{q,1}(|u|)^2 \frac{|Du|^4}{|u|^2} + (T''_{q,1}(|u|)|u| - T'_{q,1}(|u|))^2 \frac{|Du \cdot u|^4}{|u|^6} \\ &\quad - 2 \frac{q-1}{q} T'_{q,1}(|u|) (T''_{q,1}(|u|)|u| - T'_{q,1}(|u|)) \frac{|Du \cdot u|^2 |Du|^2}{|u|^4}, \end{aligned}$$

which is non-negative by Young's inequality. This finishes the proof of (iii) for the case $K = 1$. To complete the proof of the lemma it is sufficient to observe that for general $K > 0$ the coefficients a, b, c have to be replaced by $aK^{2q-2}, bK^{2q-1}, cK^{2q}$, and the conclusion then follows exactly as above. \square

4 Higher differentiability of weak solutions

In this section we start working on the solution u of the parabolic system (2.3). First, we prove an upper bound for the average of *weighted* norms of Du . This will be done in Section 4.1 and serves also to explain the general strategy to obtain such estimates. We will then extract higher regularity properties of the solution, still following the ideas given in Section 4.1. More precisely, as final result of this section, we are interested in pathwise higher integrability of the gradient Du , which will be the core of the proof of the regularity result given in Theorem 5.1.

4.1 An a priori estimate

From Definition 2.2 of a weak solution u to (2.3), no a priori information is available on the expected value of the solution. In particular, we only know that $u(\omega)$ belongs to the space $V^2(D_T, \mathbb{R}^N)$ for P -almost every $\omega \in \Omega$, but it is still possible that the average $E[\|u\|_{V^2(D_T, \mathbb{R}^N)}]$ is infinite. Even if this cannot be excluded, we can win an a priori information on the average of weighted norms of u .

The strategy for a priori estimates for deterministic elliptic or parabolic systems is simply to “test” the equation with the solution (or some modification of it), and an estimate then follows by employing the

regularity and growth properties of the system. For stochastic systems, testing with an the appropriate version is replaced by the application of an Itô's formula for Banach spaces. Then, a first pathwise estimate follows (Step 1). Since we are interested in averages, the first estimate is rewritten (Step 2) by introducing weights depending on the solution itself. With these weights we can finally take the expectation (Step 3) and end up with the desired estimate, which we now state in its precise form.

Lemma 4.1. *Let $u \in V^2(D_T, \mathbb{R}^N)$ be a weak solution to the initial boundary value problem to (2.3) under the assumptions (2.1)_{1,2}, (2.2)_{1,2} and with $u(\cdot, 0) = u_0(\cdot) \in L^2(D, \mathbb{R}^N)$. Suppose further that the smallness condition $L_H^2 < 2\kappa^{-1}(1 - (1 - \nu^2)^{1/2})$ is satisfied, and let $D_0 \subset D$ with $d_0 := \text{dist}(D_0, \partial D) > 0$. Then there holds*

$$E \left[\int_0^T e^{-\int_0^t c_0 G_0(u, f) ds} \|Du(t)\|_{L^2(D_0)}^2 dt \right] \leq c_0 (\|u_0\|_{L^2(D)}^2 + 1 + E[\|f_H\|_{L^2(D_T)}^2])$$

for a constant c_0 depending only on D, L, L_H, d_0, κ and ν , and a function $G_0(u, f)$ given by (4.3).

Proof. Step 1. A preliminary pathwise estimate. We start by multiplying the equation (2.3) with a standard cut-off function $\eta \in C^\infty(D, [0, 1])$ which satisfies $\eta \equiv 1$ on D_0 and $|D\eta| \leq c(d_0)$. Obviously, the map $\eta \triangle_{k,h} u$ has the same properties concerning integrability and measurability as u , and the Itô formula from Theorem 3.2 in Banach spaces may be applied with the Gelfand triple $W_0^{1,2}(D, \mathbb{R}^N) \subset L^2(D, \mathbb{R}^N) \subset W^{-1,2}(D, \mathbb{R}^N)$. This yields the existence of a subset $\Omega' \subset \Omega$ of full measure $P(\Omega') = 1$ and a function $u': [0, T] \rightarrow W^{1,2}(D, \mathbb{R}^N)$ which satisfies: u' is \mathcal{F}_t -adapted on $[0, T] \times \Omega'$, continuous in t for every $\omega \in \Omega'$, and $u' = u\eta$ holds for $P \times \mathcal{L}^1$ -almost all $(t, \omega) \in [0, T] \times \Omega$. Moreover, using the integration by parts formula, we have for every $\omega \in \Omega'$

$$\begin{aligned} \|u'(t)\|_{L^2(D)}^2 + 2 \int_0^t \langle D(u(s)\eta^2), A(\cdot, s, u(s), Du(s)) \rangle_{L^2(D)} ds \\ = \|u_0\eta\|_{L^2(D)}^2 + 2 \int_0^t \langle u'(s)\eta, H(\cdot, s, Du(s)) dB_s \rangle_{L^2(D)} + \int_0^t \|H(\cdot, s, Du(s))\eta\|_{L^2(D)}^2 ds \end{aligned} \quad (4.1)$$

(with the convention $|M|^2 = \sum_{i,j=1}^n M_{ij}^2$ for every $n \times n$ matrix). Next we need to estimate the second integral on the left-hand side of the previous identity, employing the assumptions (2.1). For this purpose, we first observe with (2.1)_{1,2} and Young's inequality that

$$\begin{aligned} \langle Du(s)\eta^2, A(\cdot, s, u(s), Du(s)) \rangle_{L^2(D)} \\ = \int_0^1 \langle Du(s)\eta^2, D_z A(\cdot, s, u(s), rDu(s)) Du(s) \rangle_{L^2(D)} dr + \langle Du(s)\eta^2, A(\cdot, s, u(s), 0) \rangle_{L^2(D)} \\ \geq \frac{1}{\kappa} (1 - (1 - \nu^2)^{\frac{1}{2}} - \varepsilon) \|Du(s)\eta\|_{L^2(D)}^2 - c(\varepsilon^{-1}, L) \left(\|u(s)\|_{L^{\frac{2(n+2)}{n}}(D)}^{\frac{2(n+2)}{n}} + \|f(s)\|_{L^a(D)}^a \right) \end{aligned}$$

for all $s \in (0, T)$ and every $\varepsilon > 0$. Moreover, we find

$$\begin{aligned} |\langle u(s) \otimes D\eta\eta, A(\cdot, s, u(s), Du(s)) \rangle_{L^2(D)}| \\ \leq \frac{\varepsilon}{\kappa} \|Du(s)\eta\|_{L^2(D)}^2 + c(L, \varepsilon^{-1}, \kappa, d_0) (\|u(s)\|_{L^2(D)}^2 + \|u(s)\|_{L^{\frac{2(n+2)}{n}}(D)}^{\frac{2(n+2)}{n}} + \|f(s)\|_{L^a(D)}^a). \end{aligned}$$

Next, the integrand of the last term on the right-hand side of (4.1) is bounded via (2.2)_{1,2} by

$$\|H(\cdot, s, Du(s))\eta\|_{L^2(D)}^2 \leq (L_H^2 + \frac{\varepsilon}{\kappa}) \|Du(s)\eta\|_{L^2(D)}^2 + c(\varepsilon^{-1}, \kappa) \|f_H(s)\|_{L^2(D)}^2.$$

Combining the last three inequalities (here enters the smallness assumption on L_H) with (4.1), choosing ε sufficiently small (in dependency of L_H and ν) and using Hölder's inequality, we thus end up with the

announced pathwise estimate

$$\begin{aligned}
& \|u'(t)\|_{L^2(D)}^2 + c^{-1}(L_H, \kappa, \nu) \int_0^t \|Du(s) \eta\|_{L^2(D)}^2 ds \\
& \leq \|u_0 \eta\|_{L^2(D)}^2 + 2 \int_0^t \langle u'(s) \eta, H(\cdot, s, Du(s)) dB_s \rangle_{L^2(D)} \\
& \quad + c_0(D, L, L_H, d_0, \kappa, \nu) \int_0^t \left(1 + \|u(s)\|_{L^{\frac{2(n+2)}{n}}(D)}^{\frac{2(n+2)}{n}} + \|f(s)\|_{L^a(D)}^a + \|f_H(s)\|_{L^2(D)}^2\right) ds. \tag{4.2}
\end{aligned}$$

Step 2. An improved pathwise estimate. The next step consists in getting a pathwise estimate where the bound on the right-hand side contains a deterministic part almost *independent* of the weak solution and the function f , and a stochastic part which might still depend on the solution. We start by defining

$$G_0(u, f)(s) = 1 + \|u(s)\|_{L^{\frac{2(n+2)}{n}}(D)}^{\frac{2(n+2)}{n}} + \|f(s)\|_{L^a(D)}^a \tag{4.3}$$

for $s \in (0, T)$. Obviously, G_0 belongs to $L^1(0, T)$ with probability one. Then we use a Gronwall-type argument, by applying the one-dimensional Itô-formula to $\exp(-\int_0^t c_0 G_0(u, f)(\tilde{s}) d\tilde{s}) (1 + \|u'(t)\|_{L^2(D)}^2)$, as e. g. in [25, Proof of Theorem 5.1]. Thus, we get

$$\begin{aligned}
& e^{-\int_0^t c_0 G_0(u, f)(\tilde{s}) d\tilde{s}} \|u'(t)\|_{L^2(D)}^2 + c^{-1}(L_H, \kappa, \nu) \int_0^t e^{-\int_0^s c_0 G_0(u, f)(\tilde{s}) d\tilde{s}} \|Du(s) \eta\|_{L^2(D)}^2 ds \\
& \leq \|u(0) \eta\|_{L^2(D)}^2 + 1 + 2 \int_0^t e^{-\int_0^s c_0 G_0(u, f)(\tilde{s}) d\tilde{s}} \langle u'(s) \eta, H(\cdot, s, Du(s)) dB_s \rangle_{L^2(D)} \\
& \quad + c_0 \int_0^t e^{-\int_0^s c_0 G_0(u, f)(\tilde{s}) d\tilde{s}} \|f_H(s)\|_{L^2(D)}^2 ds. \tag{4.4}
\end{aligned}$$

Note that we here have omitted a negative term which appeared on the right-hand side and the positive term containing the 1 on the left-hand side. This is the desired improved pathwise estimate. We note that u and f still appear in the function G_0 in the deterministic integral on the right-hand side, but in a way that for greater values of u or f the integral gets smaller. At the same time obviously also the exponential factor on the left-hand side will get smaller, but this allows us now to pass to Step 3.

Step 3. An estimate for the expected value with weights. Uniform estimates for the average of the weak solution (e. g. for expressions of the form $E[\|u\|]$ for some norm of u or Du) of course cannot be expected under such weak assumptions as we have supposed in the lemma. But the previous inequality (4.4) now allows us to get a weighted inequality, with no stochastic terms on the right-hand side. Since the expectation of the stochastic integral is not a priori known to vanish, we now apply a stopping time argument.

From identity (4.1) it follows that the process $\|u'(t)\|_{L^2(D)}^2$ has a continuous version in t , used in the following argument. For every $R > 0$ we introduce the random time

$$\tau_R := \inf \left\{ t \in [0, T] : \int_0^t \|u'(s)\|_{L^2(D)}^2 \|H(s, Du(s)) \eta\|_{L^2(D)}^2 ds > R \right\}$$

with $\tau_R = T$ when the set is empty. We note that $\|u'(s)\|_{L^2(D)}^2 \|H(s, Du(s)) \eta\|_{L^2(D)}^2$ is in $L^1(0, T)$ with probability one, because of the property $u \in V_0^2(D_T, \mathbb{R}^N)$ and the assumption (2.2)₂ on H . Hence, we have in particular $P(\lim_{R \rightarrow \infty} \tau_R = T) = 1$ and

$$P\left(\lim_{R \rightarrow \infty} \|u'(t \wedge \tau_R)\|_{L^2(D)}^2 = \|u'(t)\|_{L^2(D)}^2\right) = 1$$

for every $t \in [0, T]$. Now we take inequality (4.4) at time $t \wedge \tau_R$ and get

$$\begin{aligned} & e^{-\int_0^{t \wedge \tau_R} c_0 G_0(u, f)(\bar{s}) d\bar{s}} \|u'(t \wedge \tau_R)\|_{L^2(D)}^2 + c^{-1} \int_0^{t \wedge \tau_R} e^{-\int_0^s c_0 G_0(u, f)(\bar{s}) d\bar{s}} \|Du(s) \eta\|_{L^2(D)}^2 ds \\ & \leq \|u(0) \eta\|_{L^2(D)}^2 + 1 + 2 \int_0^{t \wedge \tau_R} e^{-\int_0^s c_0 G_0(u, f)(\bar{s}) d\bar{s}} \langle u'(s) \eta, H(\cdot, s, Du(s)) dB_s \rangle_{L^2(D)} \\ & \quad + c_0 \int_0^{t \wedge \tau_R} e^{-\int_0^s c_0 G_0(u, f)(\bar{s}) d\bar{s}} \|f_H(s)\|_{L^2(D)}^2 ds. \end{aligned}$$

Now we have

$$\begin{aligned} & \int_0^{t \wedge \tau_R} e^{-\int_0^s c_0 G_0(u, f)(\bar{s}) d\bar{s}} \langle u'(s) \eta, H(s, Du(s)) dB_s \rangle_{L^2(D)} \\ & = \int_0^t e^{-\int_0^s c_0 G_0(u, f)(\bar{s}) d\bar{s}} 1_{s \leq \tau_R} \langle u'(s) \eta, H(s, Du(s)) dB_s \rangle_{L^2(D)} \end{aligned}$$

and

$$\begin{aligned} & \int_0^t e^{-2 \int_0^s c_0 G_0(u, f)(\bar{s}) d\bar{s}} 1_{s \leq \tau_R} \|u'(s)\|_{L^2(D)}^2 \|H(s, Du(s)) \eta\|_{L^2(D)}^2 ds \\ & \leq \int_0^{t \wedge \tau_R} \|u'(s)\|_{L^2(D)}^2 \|H(s, Du(s)) \eta\|_{L^2(D)}^2 ds \leq R \end{aligned}$$

by definition of τ_R . Thus the stopped stochastic integral above is a martingale, hence with expected value zero. This implies

$$\begin{aligned} & E \left[e^{-\int_0^{t \wedge \tau_R} c_0 G_0(u, f)(\bar{s}) d\bar{s}} \|u'(t \wedge \tau_R)\|_{L^2(D)}^2 \right] + E \left[c^{-1} \int_0^{t \wedge \tau_R} e^{-\int_0^s c_0 G_0(u, f)(\bar{s}) d\bar{s}} \|Du(s) \eta\|_{L^2(D)}^2 ds \right] \\ & \leq \|\Delta_{k,h} u_0 \eta\|_{L^2(D)}^2 + 1 + c_0 E \left[\int_0^T \|f_H(s)\|_{L^2(D)}^2 ds \right]. \end{aligned}$$

On the left-hand-side we now apply Fatou's lemma to the first term and the monotone convergence theorem to the second one, and we get

$$\begin{aligned} & E \left[e^{-\int_0^t c_0 G_0(u, f)(\bar{s}) d\bar{s}} \|u'(t)\|_{L^2(D)}^2 \right] + E \left[c^{-1} \int_0^t e^{-\int_0^s c_0 G_0(u, f)(\bar{s}) d\bar{s}} \|Du(s) \eta\|_{L^2(D)}^2 ds \right] \\ & \leq \|u_0 \eta\|_{L^2(D)}^2 + 1 + c_0 E \left[\int_0^T \|f_H(s)\|_{L^2(D)}^2 ds \right] \end{aligned}$$

for every $t \in [0, T]$. This proves the bound claimed by the lemma. \square

4.2 Existence of second space derivatives

We next study the existence of second order space derivatives. For deterministic elliptic and parabolic partial differential equations it is a standard procedure to establish the existence of higher order derivatives by finite difference quotients methods. The basic idea in the deterministic case is the following. Once the norm of finite difference quotients of Du are kept under control *independently* of its step size, i. e. $\|\Delta_{k,h} Du\|_{L^p} \leq C$ with C independent of h and with $p \in (1, \infty)$, then the weak derivative $D_k Du$ exists and has finite norm in L^p (and as long as one is away from the boundary also the reverse is true). So uniformly bounded difference quotients of Du can heuristically be considered as second derivatives $D_k Du$. This uniform bound in turn is usually achieved by “testing the system” with appropriate modifications of the solutions (formally one might think of $\Delta_{k,-h} \Delta_{k,h} u$) and relies on the one hand on the ellipticity of the vector field A and on the other

hand on its regularity with respect to the x and u variables (we note that it seems mandatory to have at least Lipschitz-regularity in order to expect the existence of full second space derivatives).

For the stochastic perturbed system (2.3) the approach for proving the existence of higher order derivatives is still very similar, but we need some modifications due to the stochastic terms. The above strategy (with testing replaced by the use of the Itô formula in Banach spaces) applied to our stochastic system gives – after some standard, though very technical computations – a preliminary pathwise estimate for finite difference quotients of u and Du (this corresponds in some sense to Step 1 in the proof of the previous Lemma 4.1). But since this estimate still involves a stochastic integral, it is not yet possible to gain immediately any information on second order derivatives. In a second step, this pathwise estimate is rewritten (here again some Gronwall-type inequality is needed), which allows in the third step to take the expectation of a weighted version of $\|\Delta_{k,h} Du\|_{L^2}$ and to bound it independently of the stepsize h . This is still sufficient to deduce the existence of $D_k Du$ with probability one (see Section 3.2).

Given a deterministic initial condition u_0 (sufficiently regular) we now give the precise statement on the boundedness of the expectation of finite difference quotients of Du .

Lemma 4.2. *Let $u \in V^2(D_T, \mathbb{R}^N)$ be a weak solution to the initial boundary value problem to (2.3) under the assumptions (2.1), (2.2) and with $u(\cdot, 0) = u_0(\cdot) \in W^{1,2}(D, \mathbb{R}^N)$. Suppose further that the smallness condition $L_H^2 < 2\kappa^{-1}(1 - (1 - \nu^2)^{1/2})$ is satisfied, and let $D' \subset D$ with $d' := \text{dist}(D', \partial D) > 0$. Then there holds*

$$\begin{aligned} \sup_{|h| < d'} E \left[\sup_{t \in (0, T)} e^{-\int_0^t c' G'(u, f) ds} \|\Delta_{k,h} u(t)\|_{L^2(D')}^2 + \int_0^T e^{-\int_0^t c' G'(u, f) ds} \|D \Delta_{k,h} u(t)\|_{L^2(D')}^2 dt \right] \\ \leq c' (\|D_k u_0\|_{L^2(D)}^2 + 1 + E[\|f_H^{\frac{\alpha}{\alpha-2}}\|_{L^2(D_T)}^2]) \end{aligned}$$

for every $k \in \{1, \dots, n\}$, a constant c' depending only on $n, D, T, L, L_H, d', \kappa$, and ν , and a function $G'(u, f)$ given by (4.7) further below.

Proof. We here proceed similarly to the proof of Lemma 4.1, with the main difference that instead of u we now need to estimate the difference quotients $\Delta_{k,h} u$.

Step 1. We first observe that if u is a solution of (2.3), then for all $t \in [0, T]$ by definition also the following identity holds true for P -a.s.:

$$\begin{aligned} \langle \eta \Delta_{k,h} u(t) - \eta \Delta_{k,h} u_0, \varphi \rangle_{L^2(D)} &= \int_0^t \langle \eta \operatorname{div} \Delta_{k,h} A(\cdot, s, u(s), Du(s)), \varphi \rangle_{W^{-1,2}(D); W_0^{1,2}(D)} ds \\ &\quad + \int_0^t \langle \varphi, \eta \Delta_{k,h} H(\cdot, s, Du(s)) \rangle_{L^2(D)} dB_s \end{aligned}$$

for all $\varphi \in W_0^{1,2}(D, \mathbb{R}^N)$. Here $k \in \{1, \dots, n\}$ is arbitrary, $h \in \mathbb{R}$ with $|h| < d'$, and for sets $D' \subset D_0 \subset D$ we denote by $\eta \in C^\infty(D_0, [0, 1])$ a standard cut-off function satisfying $\eta \equiv 1$ on D' and $|D\eta| \leq c(d')$. Therefore, the map $\eta \Delta_{k,h} u$ has the same properties concerning integrability and measurability as u , and the Itô formula from Theorem 3.2 in Banach spaces is again applied with the Gelfand triple $W_0^{1,2}(D, \mathbb{R}^N) \subset L^2(D, \mathbb{R}^N) \subset W^{-1,2}(D, \mathbb{R}^N)$. We hence get a subset $\Omega' \subset \Omega$ of full measure $P(\Omega') = 1$ and a function $u'_k : [0, T] \rightarrow W^{1,2}(D, \mathbb{R}^N)$ with the following properties: u'_k is \mathcal{F}_t -adapted on $[0, T] \times \Omega'$, continuous in t for every $\omega \in \Omega'$, and satisfies $u'_k = \Delta_{k,h} u \eta$ for $P \times \mathcal{L}^1$ -almost all $(t, \omega) \in [0, T] \times \Omega$. Moreover, for every $\omega \in \Omega'$ we have

$$\begin{aligned} \|u'_k(t)\|_{L^2(D)}^2 + 2 \int_0^t \langle D(\Delta_{k,h} u(s) \eta^2), \Delta_{k,h} A(\cdot, s, u(s), Du(s)) \rangle_{L^2(D)} ds \\ = \|\Delta_{k,h} u_0 \eta\|_{L^2(D)}^2 + 2 \int_0^t \langle u'_k(s) \eta, \Delta_{k,h} H(\cdot, s, Du(s)) \rangle_{L^2(D)} dB_s \\ + \int_0^t \|\Delta_{k,h} H(\cdot, s, Du(s)) \eta\|_{L^2(D)}^2 ds. \end{aligned} \tag{4.5}$$

Our first aim is to deduce a pathwise estimate for finite differences of u and Du , respectively. To this end we first study in detail the second term on the left-hand side. For almost every $t \in [0, T]$ we decompose the finite difference quotient applied on $A(x, t, u, Du)$ as follows

$$\begin{aligned}
& \Delta_{k,h} A(x, t, u(x), Du(x)) \\
&= h [A(x + he_k, t, u(x + he_k), Du(x + he_k)) - A(x + he_k, t, u(x + he_k), Du(x))] \\
&\quad + h [A(x + he_k, t, u(x + he_k), Du(x)) - A(x + e_k, t, u(x), Du(x))] \\
&\quad + h [A(x + he_k, t, u(x), Du(x)) - A(x, t, u(x), Du(x))] \\
&= \int_0^1 D_z A(x + he_k, t, u(x + he_k), Du(x) + r h \Delta_{k,h} Du(x)) dr \Delta_{k,h} Du(x) \\
&\quad + \int_0^1 D_u A(x + he_k, t, u(x) + r h \Delta_{k,h} u(x), Du(x)) dr \Delta_{k,h} u(x) \\
&\quad + \int_0^1 D_x A(x + r h e_k, t, u(x), Du(x)) dr \\
&=: \mathcal{A}(h) + \mathcal{B}(h) + \mathcal{C}(h)
\end{aligned} \tag{4.6}$$

with the obvious abbreviations. Using the assumptions (2.1), Hölder's and Young's inequality, we now estimate the different terms arising from this decomposition in equation (4.5) on time slices $t \in (0, T)$ (on such slices we omit the notion of t). We first find for almost every $t \in (0, T)$ and every $\varepsilon > 0$

$$\begin{aligned}
& \langle D(\Delta_{k,h} u \eta^2), \mathcal{A}(h) \rangle_{L^2(D)} \\
&= \kappa^{-1} \langle D(\Delta_{k,h} u \eta^2), D\Delta_{k,h} u \rangle_{L^2(D)} - \kappa^{-1} \langle D(\Delta_{k,h} u \eta^2), D\Delta_{k,h} u - \kappa \mathcal{A}(h) \rangle_{L^2(D)} \\
&\geq \kappa^{-1} \|D\Delta_{k,h} u \eta\|_{L^2(D)}^2 - 2\varepsilon \|D\Delta_{k,h} u \eta\|_{L^2(D)}^2 - c(\kappa, \varepsilon) \|\Delta_{k,h} u D\eta\|_{L^2(D)}^2 \\
&\quad - \kappa^{-1} (1 - \nu^2)^{\frac{1}{2}} \|D\Delta_{k,h} u \eta\|_{L^2(D)}^2 \\
&\geq (\kappa^{-1} (1 - (1 - \nu^2)^{\frac{1}{2}}) - 2\varepsilon) \|D\Delta_{k,h} u \eta\|_{L^2(D)}^2 - c(\kappa, \varepsilon, \|D\eta\|_{L^\infty(D)}) \|D_k u\|_{L^2(D)}^2
\end{aligned}$$

where in the last line we have used the fact that the norm of the finite difference quotient of a compactly supported function is always bounded by the norm of the partial derivative (provided that the stepsize is sufficiently small). This lower bound will be crucial (and can be understood as some ellipticity of the vector field A up to lower order terms). We next observe with the Sobolev-Poincaré embedding (applied on every time-slice)

$$\begin{aligned}
& |\langle D(\Delta_{k,h} u \eta^2), \mathcal{B}(h) \rangle_{L^2(D)}| \\
&\leq L (\|D\Delta_{k,h} u \eta\|_{L^2(D)} + 2\|\Delta_{k,h} u D\eta\|_{L^2(D)}) \|\Delta_{k,h} u \eta\|_{L^{\frac{2n}{n-2}}(D)}^\theta \\
&\quad \times \|\Delta_{k,h} u \eta\|_{L^2(D)}^{1-\theta} (\|Du\|_{L^{\frac{2}{(n+2)\theta}}(D)}^{\frac{2}{n+2}} + \|u\|_{L^{\frac{2}{\theta}}(D)}^{\frac{2}{n}} + \|f\|_{L^{\frac{n}{\theta}}(D)}) \\
&\leq (\|D\Delta_{k,h} u \eta\|_{L^2(D)}^{1+\theta} + \|\Delta_{k,h} u D\eta\|_{L^2(D)}^{1+\theta}) \\
&\quad \times c(n, D, L) \|\Delta_{k,h} u \eta\|_{L^2(D)}^{1-\theta} (\|Du\|_{L^{\frac{2}{(n+2)\theta}}(D)}^{\frac{2}{n+2}} + \|u\|_{L^{\frac{2}{\theta}}(D)}^{\frac{2}{n}} + \|f\|_{L^{\frac{n}{\theta}}(D)})
\end{aligned}$$

for every $\theta \in (0, 1)$; in the two-dimensional case $n = 2$, $\frac{2n}{n-2}$ shall be interpreted as any arbitrary number greater than 2. We note that we have omitted the step of passing from the shifted to the original domain in the first inequality and we have applied $cd^\theta + c^\theta d \leq c^{1+\theta} + d^{1+\theta}$ for all $c, d \geq 0$ to get from the first to the second inequality. To estimate further we choose $\theta = \frac{n}{n+2}$, according to the integrability assumptions of u (using the embedding given in (2.4)), Du and f . Thus, Young's inequality gives

$$\begin{aligned}
& |\langle D(\Delta_{k,h} u \eta^2), \mathcal{B}(h) \rangle_{L^2(D)}| \\
&\leq \varepsilon \|D\Delta_{k,h} u \eta\|_{L^2(D)}^2 + c(n, D, L, \|D\eta\|_{L^\infty(D)}, \varepsilon) (\|\Delta_{k,h} u \eta\|_{L^2(D)}^2 + 1)
\end{aligned}$$

$$\times \left(\|Du\|_{L^2(D)}^2 + \|u\|_{L^{\frac{2(n+2)}{n}}(D)}^{\frac{2(n+2)}{n}} + \|f\|_{L^{n+2}(D)}^{n+2} \right).$$

Finally, using Young's inequality and standard properties of finite difference quotients, we estimate the last term involving $\mathcal{C}(h)$ by

$$\begin{aligned} & \left| \langle D(\Delta_{k,h} u \eta^2), \mathcal{C}(h) \rangle_{L^2(D)} \right| \\ & \leq L \left(\|D\Delta_{k,h} u \eta\|_{L^2(D)} + 2\|\Delta_{k,h} u D\eta\|_{L^2(D)} \right) \left(\|Du\|_{L^2(D)} + \|u\|_{L^{\frac{2(n+2)}{n}}(D)}^{\frac{n+2}{n}} + \|f\|_{L^4(D)}^2 \right) \\ & \leq \varepsilon \|D\Delta_{k,h} u \eta\|_{L^2(D)}^2 + c(D, L, \|D\eta\|_{L^\infty(D)}, \varepsilon) \left(1 + \|Du\|_{L^2(D)}^2 + \|u\|_{L^{\frac{2(n+2)}{n}}(D)}^{\frac{2(n+2)}{n}} + \|f\|_{L^{n+2}(D)}^{n+2} \right). \end{aligned}$$

Now we have estimated all terms coming from the integral involving $\Delta_{k,h} A(x, t, u, Du)$. Next we study the last integral in equation (4.5). Employing the properties (2.2), we find

$$\|\Delta_{k,h} H(s, Du(s)) \eta\|_{L^2(D)} \leq \|f_H^{\frac{a}{a-2}}(s) + |Du(s)| \eta\|_{L^2(D)} + L_H \|D\Delta_{k,h} u(s) \eta\|_{L^2(D)}.$$

Before summarizing the previous estimates for the single terms, we introduce, for ease of notation, the function

$$G'(u, f)(s) := 1 + \|Du(s)\|_{L^2(D)}^2 + \|u(s)\|_{L^{\frac{2(n+2)}{n}}(D)}^{\frac{2(n+2)}{n}} + \|f(s)\|_{L^a(D)}^a \quad (4.7)$$

with $s \in (0, T)$, which belongs to $L^1(0, T)$ almost surely. Note that by definition we have $G'(u, f) \geq G_0(u, f)$ with $G_0(u, f)$ denoting the function introduced in Lemma 4.1. Combining the latter estimates with the decomposition given in (4.6), using standard properties for finite difference quotients and choosing $\varepsilon = \varepsilon(\kappa, \nu, L_H)$ sufficiently small, we hence infer from (4.5) that for every $\omega \in \Omega'$ there holds

$$\begin{aligned} & \|u'_k(t)\|_{L^2(D)}^2 + c^{-1}(L_H, \kappa, \nu) \int_0^t \|D\Delta_{k,h} u(s) \eta\|_{L^2(D)}^2 ds \\ & \leq \|\Delta_{k,h} u(0) \eta\|_{L^2(D)}^2 + c' \int_0^t (\|\Delta_{k,h} u(s) \eta\|_{L^2(D)}^2 + 1) G'(u, f)(s) ds \\ & \quad + 2 \int_0^t \langle u'_k(s) \eta, \Delta_{k,h} H(\cdot, s, Du(s)) dB_s \rangle_{L^2(D)} + 2 \int_0^t \|f_H^{\frac{a}{a-2}}(s)\|_{L^2(D)}^2 ds, \end{aligned} \quad (4.8)$$

and the constant c' depends only on n, D, L, L_H, d', κ and ν . Here we assume $c' \geq c_0$ with c_0 denoting the constant given in Lemma 4.1. This is the preliminary pathwise estimate on the finite difference quotients (which however involves the stochastic integral) and concludes the Step 1.

Step 2. Before passing to the expectation value as described in the beginning we still need the Gronwall-type argument, similarly as in the proof of Lemma 4.1. However, we here observe that the second integral on the right-hand side is in general not known to be finite, but the first factor of the integrand “almost” happens to appear in the sum of its left-hand side (in the sense that u'_k differs from $\Delta_{k,h} u \eta$ only on a negligible set). Hence, to get rid of this possibly uncontrollable term we apply the one-dimensional Itô-formula (recalling $a > n + 2$), and we obtain

$$\begin{aligned} & e^{-\int_0^t c' G'(u, f)(\bar{s}) d\bar{s}} \|u'_k(t)\|_{L^2(D)}^2 + c^{-1} \int_0^t e^{-\int_0^s c' G'(u, f)(\bar{s}) d\bar{s}} \|D\Delta_{k,h} u(s) \eta\|_{L^2(D)}^2 ds \\ & \leq c' \int_0^t e^{-\int_0^s c' G'(u, f)(\bar{s}) d\bar{s}} G'(u, f) (\|\Delta_{k,h} u(s) \eta\|_{L^2(D)}^2 - \|u'_k(s)\|_{L^2(D)}^2) ds \\ & \quad + \|\Delta_{k,h} u(0) \eta\|_{L^2(D)}^2 + 1 + 2 \int_0^t e^{-\int_0^s c' G'(u, f)(\bar{s}) d\bar{s}} \langle u'_k(s) \eta, \Delta_{k,h} H(\cdot, s, Du(s)) dB_s \rangle_{L^2(D)} \\ & \quad + 2 \int_0^t e^{-\int_0^s c' G'(u, f)(\bar{s}) d\bar{s}} \|f_H^{\frac{a}{a-2}}(s)\|_{L^2(D)}^2 ds. \end{aligned} \quad (4.9)$$

Here, we again note that the average of the first integral on the right-hand side vanishes, due to the fact that u'_k and $\Delta_{k,h}u\eta$ coincide on $[0, T] \times \Omega$ except for a set of $\mathcal{L}^1 \times P$ -measure zero.

Step 3. We now derive the estimate for the average with weights. In contrast to Lemma 4.1, we now derive in a first step an estimate for the (weighted) average of the $L^2(L^2)$ -norm of $D\Delta_{k,h}u$ (which proceeds in exactly the same way as before). Then, we use this estimate to get also an upper bound for the (weighted) average of the $L^\infty(L^2)$ -norm of $\Delta_{k,h}u$.

We first note that identity (4.5) implies that the process $\|u'_k(t)\|_{L^2(D)}^2$ has a continuous version in t . Now, for every $R > 0$, we introduce the random time

$$\tau_R := \inf \left\{ t \in [0, T] : \int_0^t \|u'_k(s)\|_{L^2(D)}^2 \|\Delta_{k,h}H(s, Du(s))\eta\|_{L^2(D)}^2 ds > R \right\}$$

with $\tau_R = T$ when the set is empty. Notice that

$$\int_0^T \|u'_k(s)\|_{L^2(D)}^2 \|\Delta_{k,h}H(s, Du(s))\eta\|_{L^2(D)}^2 ds < \infty$$

with probability one, because of the property $u \in V_0^2(D_T, \mathbb{R}^N)$ and the assumption (2.2)₂ on H . Hence, when $R \rightarrow \infty$, τ_R is eventually equal to T , with probability one. In particular, $P(\lim_{R \rightarrow \infty} \tau_R = T) = 1$ and for every $t \in [0, T]$ we have

$$P\left(\lim_{R \rightarrow \infty} \|u'_k(t \wedge \tau_R)\|_{L^2(D)}^2 = \|u'_k(t)\|_{L^2(D)}^2\right) = 1.$$

Step 3a. We compute inequality (4.9) at time $t \wedge \tau_R$ and get

$$\begin{aligned} & e^{-\int_0^{t \wedge \tau_R} c' G'(u, f)(\bar{s}) d\bar{s}} \|u'_k(t \wedge \tau_R)\|_{L^2(D)}^2 + c^{-1} \int_0^{t \wedge \tau_R} e^{-\int_0^s c' G'(u, f)(\bar{s}) d\bar{s}} \|D\Delta_{k,h}u(s)\eta\|_{L^2(D)}^2 ds \\ & \leq c' \int_0^{t \wedge \tau_R} e^{-\int_0^s c' G'(u, f)(\bar{s}) d\bar{s}} G'(u, f) (\|\Delta_{k,h}u(s)\eta\|_{L^2(D)}^2 - \|u'_k(s)\|_{L^2(D)}^2) ds \\ & \quad + \|\Delta_{k,h}u_0\eta\|_{L^2(D)}^2 + 1 + 2 \int_0^{t \wedge \tau_R} e^{-\int_0^s c' G'(u, f)(\bar{s}) d\bar{s}} \langle u'_k(s)\eta, \Delta_{k,h}H(s, Du(s)) dB_s \rangle_{L^2(D)} ds \\ & \quad + 2 \int_0^{t \wedge \tau_R} e^{-\int_0^s c' G'(u, f)(\bar{s}) d\bar{s}} \|f_H^{\frac{a}{a-2}}(s)\|_{L^2(D)}^2 ds. \end{aligned}$$

Now we have

$$\begin{aligned} & \int_0^{t \wedge \tau_R} e^{-\int_0^s c' G'(u, f)(\bar{s}) d\bar{s}} \langle u'_k(s)\eta, \Delta_{k,h}H(s, Du(s)) dB_s \rangle_{L^2(D)} \\ & = \int_0^t e^{-\int_0^s c' G'(u, f)(\bar{s}) d\bar{s}} 1_{s \leq \tau_R} \langle u'_k(s)\eta, \Delta_{k,h}H(s, Du(s)) dB_s \rangle_{L^2(D)} \end{aligned}$$

and

$$\begin{aligned} & \int_0^t e^{-2 \int_0^s c' G'(u, f)(\bar{s}) d\bar{s}} 1_{s \leq \tau_R} \|u'_k(s)\|_{L^2(D)}^2 \|\Delta_{k,h}H(s, Du(s))\eta\|_{L^2(D)}^2 ds \\ & \leq \int_0^{t \wedge \tau_R} \|u'_k(s)\|_{L^2(D)}^2 \|\Delta_{k,h}H(s, Du(s))\eta\|_{L^2(D)}^2 ds \leq R \end{aligned}$$

by definition of τ_R . Thus, the stopped stochastic integral above is a martingale, hence with expected value zero. This implies (using that u'_k equals $\Delta_{k,h}u\eta$ outside a set of $\mathcal{L}^1 \times P$ -measure zero)

$$\begin{aligned} & E\left[e^{-\int_0^{t \wedge \tau_R} c' G'(u, f)(\bar{s}) d\bar{s}} \|u'_k(t \wedge \tau_R)\|_{L^2(D)}^2\right] + c^{-1} E\left[\int_0^{t \wedge \tau_R} e^{-\int_0^s c' G'(u, f)(\bar{s}) d\bar{s}} \|D\Delta_{k,h}u(s)\eta\|_{L^2(D)}^2 ds\right] \\ & \leq \|\Delta_{k,h}u_0\eta\|_{L^2(D)}^2 + 1 + 2 E\left[\int_0^T \|f_H^{\frac{a}{a-2}}(s)\|_{L^2(D)}^2 ds\right]. \end{aligned}$$

We apply Fatou's lemma to the first term and monotone convergence theorem to the second one on the left-hand-side, and we get

$$\begin{aligned} E \left[e^{-\int_0^t c' G'(u,f)(\bar{s}) d\bar{s}} \|u'_k(t)\|_{L^2(D)}^2 \right] + c^{-1} E \left[\int_0^t e^{-\int_0^s c' G'(u,f)(\bar{s}) d\bar{s}} \|D \Delta_{k,h} u(s) \eta\|_{L^2(D)}^2 ds \right] \\ \leq \|\Delta_{k,h} u_0 \eta\|_{L^2(D)}^2 + 1 + 2 E \left[\int_0^T \|f_H^{\frac{a}{a-2}}(s)\|_{L^2(D)}^2 ds \right] \end{aligned}$$

for every $t \in [0, T]$. This proves one of the two bounds claimed by the Lemma.

Step 3b. It is almost our final estimate except that we need the supremum in time inside the first expected value, and thus we have to repeat the previous computations by means of martingale inequalities. The previous estimate (as well as the a priori estimate from Lemma 4.1) will be used in the next one; we found it convenient to proceed in two steps. From the stopped inequality above we have

$$\begin{aligned} E \left[\sup_{t \in [0, T]} e^{-\int_0^{t \wedge \tau_R} c' G'(u,f)(\bar{s}) d\bar{s}} \|u'_k(t \wedge \tau_R)\|_{L^2(D)}^2 \right] \\ \leq \|\Delta_{k,h} u_0 \eta\|_{L^2(D)}^2 + 1 + 2 E \left[\int_0^T \|f_H^{\frac{a}{a-2}}(s)\|_{L^2(D)}^2 ds \right] \\ + 2 E \left[\sup_{t \in [0, T]} \left| \int_0^t e^{-\int_0^s c' G'(u,f)(\bar{s}) d\bar{s}} 1_{s \leq \tau_R} \langle u'_k(s) \eta, \Delta_{k,h} H(s, Du(s)) dB_s \rangle_{L^2(D)} \right| \right]. \end{aligned}$$

We apply again Fatou's lemma to the expected value on the left-hand-side. The last term on the right-hand-side is estimated, by means of the Burkholder-Davis-Gundy inequality, by

$$\begin{aligned} C E \left[\left(\int_0^T e^{-2 \int_0^s c' G'(u,f)(\bar{s}) d\bar{s}} 1_{s \leq \tau_R} \|u'_k(s)\|_{L^2(D)}^2 \|\Delta_{k,h} H(s, Du(s)) \eta\|_{L^2(D)}^2 ds \right)^{1/2} \right] \\ = C E \left[\left(\int_0^T e^{-2 \int_0^{s \wedge \tau_R} c' G'(u,f)(\bar{s}) d\bar{s}} 1_{s \leq \tau_R} \|u'_k(s \wedge \tau_R)\|_{L^2(D)}^2 \|\Delta_{k,h} H(s, Du(s)) \eta\|_{L^2(D)}^2 ds \right)^{1/2} \right] \\ \leq C E \left[I_1^{1/2} I_2^{1/2} \right] \leq \frac{1}{2} E[I_1] + \frac{C^2}{2} E[I_2] \end{aligned}$$

where

$$\begin{aligned} I_1 &= \sup_{t \in [0, T]} e^{-\int_0^{t \wedge \tau_R} c' G'(u,f)(\bar{s}) d\bar{s}} \|u'_k(t \wedge \tau_R)\|_{L^2(D)}^2, \\ I_2 &= \int_0^T e^{-\int_0^s c' G'(u,f)(\bar{s}) d\bar{s}} \|\Delta_{k,h} H(s, Du(s)) \eta\|_{L^2(D)}^2 ds. \end{aligned}$$

Hence, we have proved

$$\frac{1}{2} E[I_1] \leq \|\Delta_{k,h} u_0 \eta\|_{L^2(D)}^2 + 1 + 2 E \left[\int_0^T \|f_H^{\frac{a}{a-2}}(s)\|_{L^2(D)}^2 ds \right] + \frac{C^2}{2} E[I_2].$$

From the estimate above for $\|\Delta_{k,h} H(s, Du(s)) \eta\|_{L^2(D)}^2$ and the estimate of Step 3a, we know that $E[I_2]$ is bounded from above via

$$\begin{aligned} E[I_2] &\leq 4 E \left[\int_0^T e^{-\int_0^s c' G'(u,f)(\bar{s}) d\bar{s}} \|Du(s) \eta\|_{L^2(D)}^2 ds \right] \\ &\quad + c \left(\|\Delta_{k,h} u_0 \eta\|_{L^2(D)}^2 + 1 + E \left[\int_0^T \|f_H^{\frac{a}{a-2}}(s)\|_{L^2(D)}^2 ds \right] \right). \end{aligned}$$

With a suitable choice of D_0 (in dependency of D and D') and keeping in mind $c' G'(u, f) \geq c_0 G_0(u, f)$ by construction, the first average on the right-hand side of the last inequality is bounded due to Lemma 4.1. Hence we get the bound for $E[I_1]$ as asserted in the statement of the lemma. Since h was arbitrary and $\eta = 1$ on D' , the proof is complete. \square

Remark 4.3. For SPDEs having first space derivatives of the solution in the coefficient of the noise, the most general condition for existence of solutions in L^2 , which becomes also a condition for an improvement of $W^{k,2}$ -regularity, is more precise than just the control on the Lipschitz constant of H expressed by the statement of Lemma 4.2; see [22, 16]. However, when we go to $W^{k,p}$ -regularity with $p > 2$, the computations are too involved and the algebraic simplicity of the condition of [22, 16] seems to be lost. For this reason we have simplified the estimate also for $p = 2$.

Applying Theorem 3.6 with $(p, q) = (2, 2)$ and $(p, q) = (2, \infty)$ and summing over $k \in \{1, \dots, n\}$ we then infer from the previous lemma that second order spatial derivatives of u exist almost surely. We should note that this result does not extend up to the boundary of D since the constant c' blows up for $\text{dist}(D', D) \searrow 0$, but the result holds on any fixed subset $D' \Subset D$.

Corollary 4.4. Let $u \in V^2(D_T, \mathbb{R}^N)$ be a weak solution under the assumptions of the Lemma 4.2. Then there holds $Du \in V^2(D'_T, \mathbb{R}^N)$ with probability one, and

$$\begin{aligned} E \left[\sup_{t \in (0, T)} e^{-\int_0^t c' G'(u, f) ds} \|Du\|_{L^2(D')}^2 + \int_0^T e^{-\int_0^t c' G'(u, f) ds} \|D^2 u\|_{L^2(D')}^2 dt \right] \\ \leq c' \left(\|Du_0\|_{L^2(D)}^2 + 1 + E[\|f_H^{\frac{a}{a-2}}\|_{L^2(D_T)}^2] \right) \end{aligned}$$

for the constant c' from Lemma 4.2. Moreover, we have for all $k \in \{1, \dots, n\}$

$$e^{-\frac{1}{2} \int_0^t c' G'(u, f) ds} \triangle_{k,h} Du \rightarrow e^{-\frac{1}{2} \int_0^t c' G'(u, f) ds} D_k Du \quad \text{weakly in } L^2(D'_T \times \Omega).$$

4.3 Iteration

In the next step we want to iterate the procedure from the previous section, in a way such that we do not only know the spacial gradient Du to belong to the space V^2 with probability one, but that we get this result also for certain powers of Du . For convenience we introduce the function

$$W_q: \mathbb{R}^k \rightarrow \mathbb{R} \quad \text{defined by} \quad W_q(\xi) := |\xi|^q$$

for every $q \geq 0$. We start by briefly describing the strategy how this regularity improvement is achieved. First we observe from the results in Section 4.2 that there exists a subset of Ω of full measure on which Du belongs to $V_{\text{loc}}^2(D_T)$, hence we can now take advantage of higher integrability properties for u and Du . This shall be done with the following (but formal) iteration scheme:

$$\begin{aligned} W_{q_j}(Du) \in V^2 &\longrightarrow Du \in L^{2q_j \frac{n+2}{n}} \text{ and } u \in L^{2q_j \left(\frac{n+2}{n}\right)^2} \cap L^\infty(L^{2q_j \frac{n}{n-2}}) \\ &\longrightarrow Du A(x, t, u, Du) \in L^{\min\{q_j \frac{(n+2)^2}{n}, a\}} \text{ and } D_x A(x, t, u, Du) \in L^{\min\{2q_j \frac{n+2}{n}, \frac{a}{2}\}} \\ &\longrightarrow W_{q_{j+1}}(Du) \in V^2 \end{aligned}$$

for a sequence $\{q_j\}_{j \in \mathbb{N}}$ of numbers $q_j \geq 1$ for all $j \in \mathbb{N}$. The first implication indeed follows from the Sobolev's embedding for the space $V^{2,p}$, the second one from the growth conditions on the vector field A , the third one from the iteration (and a convergence result concerning finite difference quotients). After a finite number of steps we then arrive at a final (maximal) higher integrability exponent, which essentially reflects how close the vector field A is to the Laplace system. This should be understood in the following sense: the closer ν is to one (note that $\nu = 1$ corresponds to the case $A(x, t, u, z) = z$ plus potential lower order terms), the more integrability for Du can be gained in the iteration and the better will be the final regularity properties of u . Finally, we note that in every step of the iteration we will have to reduce the radius of the parabolic cylinder and we will also have to restrict ourselves to smaller subsets of Ω . Nevertheless, the higher integrability results will always be true on sets of probability one.

We now start with some preliminary remarks and consider again the equation (2.3)

$$du = \text{div } A(x, t, u, Du) dt + H(x, t, Du) dB_t$$

in D_T . We observe that $\operatorname{div} A(x, t, u, Du)$ is well defined in view of the regularity assumptions (2.1) and the existence of second order spatial derivatives, see Corollary 4.4. More precisely, it is easy to check that for every weak solution $u \in V_0^2(D_T, \mathbb{R}^N)$ we have: $\operatorname{div} A(x, t, u, Du) \in L_{\operatorname{loc}}^2(D', \mathbb{R}^N)$ with probability one, and the equation above holds for \mathcal{L}^n -almost every $x \in D'$ for $\mathcal{L}^1 \times P$ -almost all $(t, \omega) \in (0, T) \times \Omega$. Hence, we can now work immediately with this equation without passing to its weak formulation.

In the next lemma we will provide the main step of the iteration argument:

Lemma 4.5. *Let $u \in V^2(D_T, \mathbb{R}^N)$ be a weak solution to the initial boundary value problem to (2.3) under the assumptions (2.1), (2.2) and with initial values $u(\cdot, 0) = u_0(\cdot) \in W^{1,2}(D, \mathbb{R}^N)$, and assume that*

$$E \left[\|Y_p^p W_p(Du)\|_{L^{\frac{2p+2}{n}}(D'_T)}^{\frac{2}{p}} \right] \leq C_p < \infty$$

for some $p \geq 1$, a set $D' \subset D$, and $Y_p: [0, T] \times \Omega \rightarrow (0, 1]$ given by $Y_p(t, \omega) = \exp(-\int_0^t G_p(s, \omega) ds)$ for some function G_p which is in $L^1(0, T)$ with probability one. Let $D'' \subset D'$ with $d := \operatorname{dist}(D'', \partial D') > 0$. Then for every number $q \geq 1$ satisfying

$$q \leq \min \left\{ p \frac{n+2}{n}, 1 + p \frac{n+2}{n} \frac{a-4}{a}, \frac{a-2}{2} \right\} \quad \text{and} \quad L_H^2 < \frac{1}{\kappa(q - \frac{1}{2})} \left(\left[1 - \left(\frac{q-1}{q} \right)^2 \right]^{\frac{1}{2}} - [1 - \nu^2]^{\frac{1}{2}} \right), \quad (4.10)$$

all initial values $u_0 \in W^{1,2q}(D, \mathbb{R}^N)$, and every $k \in \{1, \dots, n\}$ there holds

$$\begin{aligned} \sup_{|h| < d} E \left[\left(\sup_{t \in (0, T)} \|Y_q^q W_q(\Delta_{k,h} u)\|_{L^2(D'')}^2 + \int_0^T \|Y_q^q DW_q(\Delta_{k,h} u)\|_{L^2(D'')}^2 dt \right)^{\frac{1}{q}} \right] \\ \leq c \left(\|W_q(D_k u_0)\|_{L^2(D)}^2 + 1 + E[\|f_H(s)\|_{L^a(D_T)}^a] \right)^{\frac{1}{q}}, \end{aligned}$$

for $Y_q: [0, T] \times \Omega \rightarrow (0, 1]$ given by $Y_q(t, \omega) = \exp(-\int_0^t G_q(s, \omega) ds)$ for some function G_q which is in $L^1(0, T)$ with probability one, and a constant c depending only on $n, p, D, T, L, L_H, d, \kappa, \nu$, and C_p .

Remarks 4.6. In the case of additive noise (with $L_H = 0$) the second condition (4.10) for the restriction on the integrability exponent q reduces to the inequality $q < \frac{1}{1-\nu}$. For multiplicative noise instead, the right-hand side in the second inequality (4.10) is decreasing in q (note that for $q = 1$ it just reproduces the condition required in Lemma 4.2) and allows the following interpretation. Obviously, the previous restriction $q < \frac{1}{1-\nu}$ for additive noise remains valid, and in fact the more multiplicative noise is considered (in the sense that L_H should not be too small), the smaller will be the maximal integrability exponent which still satisfies both inequalities in (4.10). For this reason multiplicative noise might destroy some regularity in form of integrability of the gradient Du .

Moreover, we comment on the scaling of the hypothesis and the assertion with respect to u and u_0 , respectively, in order to avoid confusion. In view of the definition of W_p it is easy to see that Lemma 4.5 is stated in a way such that an weighted average of a quadratic quantity in Du gives an information about the weighted average of a quadratic quantity in $\Delta_{k,h} u$. In this sense, the scaling is the natural one.

Proof. We now follow the line of arguments from the proof of Lemma 4.2 (and of Lemma 4.1), but this time we will estimate powers of the difference quotients $\Delta_{k,h} u$.

Step 1. We consider $k \in \{1, \dots, n\}$ arbitrary, $h \in \mathbb{R}$ with $|h| < d$, and $\eta \in C^\infty(D', [0, 1])$ a standard cut-off function satisfying $\eta \equiv 1$ on $D'' \Subset D'$ and $|\partial \eta| \leq c(d)$. We first observe that, by the integrability assumption on Du and the integrability assumption on G_p (which implies strict positivity of $\inf_{t \in [0, T]} Y_p$ for P -almost every ω), with probability one we have

$$u \in L_{\operatorname{loc}}^{2p(\frac{n+2}{n})^2}(D'_T, \mathbb{R}^N) \cap L_{\operatorname{loc}}^\infty(0, T; L^{2p\frac{n}{n-2}}(D', \mathbb{R}^N)).$$

Furthermore, due to the restriction on q , it is guaranteed that $W_q(Du)$ belongs to L^2 locally on D'_T with probability one. For almost every (fixed) $x \in D$ we first consider finite differences in direction e_k and stepsize h of the differential equation (2.3), i.e.

$$d \eta^{\frac{1}{q}} \Delta_{k,h} u(x, t) = \eta^{\frac{1}{q}} \operatorname{div} \Delta_{k,h} A(x, t, u, Du) dt + \eta^{\frac{1}{q}} \Delta_{k,h} H(x, t, Du) dB_t$$

in $(0, T)$ for $q \geq 1$. We next introduce (because of technical reasons) for $K > 0$ the approximating function $T_{q,K}(\cdot)$ of class C^2 according to Lemma 3.11, and we recall that $T_{q,K}$ satisfies in particular the polynomial growth conditions $T_{q,K}(t) = t^{2q}$ for all $t \leq K$ and $T_{q,K}(t) \leq c(q)K^{2q-2}t^2$ for all $t \in \mathbb{R}$. Employing the one-dimensional Itô formula (note that $\operatorname{div} A(x, t, u, Du)$ is as a composition of \mathcal{F}_t -adapted functions again \mathcal{F}_t -adapted) from Theorem 3.1, applied with $g(t, u(x, t)) = \eta^2 T_{q,K}(|\Delta_{k,h}u(x, t)|)$, we obtain the identity

$$\begin{aligned} & d(\eta^2 T_{q,K}(|\Delta_{k,h}u(x, t)|)) \\ &= \eta^2 T'_{q,K}(|\Delta_{k,h}u(x, t)|) |\Delta_{k,h}u(x, t)|^{-1} \langle \Delta_{k,h}u(x, t), \operatorname{div} \Delta_{k,h}A(x, t, u, Du) \rangle_{\mathbb{R}^N} dt \\ &+ \frac{1}{2} \eta^2 [T''_{q,K}(|\Delta_{k,h}u(x, t)|) |\Delta_{k,h}u(x, t)| - T'_{q,K}(|\Delta_{k,h}u(x, t)|)] \\ &\quad \times |\Delta_{k,h}u(x, t)|^{-3} |\langle \Delta_{k,h}u(x, t), \Delta_{k,h}H(x, t, Du) \rangle|^2 dt \\ &+ \frac{1}{2} \eta^2 T'_{q,K}(|\Delta_{k,h}u(x, t)|) |\Delta_{k,h}u(x, t)|^{-1} |\Delta_{k,h}H(x, t, Du)|^2 dt \\ &+ \eta^2 T'_{q,K}(|\Delta_{k,h}u(x, t)|) |\Delta_{k,h}u(x, t)|^{-1} \langle \Delta_{k,h}u(x, t), \Delta_{k,h}H(x, t, Du) dB_t \rangle_{\mathbb{R}^N}. \end{aligned}$$

In order to prove the assertion of the lemma, we start with a simple observation concerning the terms involving $\Delta_{k,h}H(x, t, Du)$. Taking into account the properties of the function $T_{q,K}$, see Lemma 3.11, we estimate

$$\begin{aligned} & [T''_{q,K}(|\Delta_{k,h}u(x, t)|) |\Delta_{k,h}u(x, t)| - T'_{q,K}(|\Delta_{k,h}u(x, t)|)] |\Delta_{k,h}u(x, t)|^{-3} |\langle \Delta_{k,h}u(x, t), \Delta_{k,h}H(x, t, Du) \rangle|^2 \\ & \leq 2(q-1) T'_{q,K}(|\Delta_{k,h}u(x, t)|) |\Delta_{k,h}u(x, t)|^{-1} |\Delta_{k,h}H(x, t, Du)|^2. \end{aligned}$$

We next introduce the abbreviation

$$V(\xi) := T'_{q,K}(|\xi|) |\xi|^{-1} \xi$$

for all $\xi \in \mathbb{R}^N$, and we note $|V(\xi)| = T'_{q,K}(|\xi|)$. Now we integrate over $x \in D$, and then we apply Fubini which due to the truncation procedure is always allowed, see Lemma 3.11 ii). Applying the integration by parts formula, we hence obtain

$$\begin{aligned} & \|(T_{q,K} |\Delta_{k,h}u(t)|)^{\frac{1}{2}} \eta\|_{L^2(D)}^2 + \int_0^t \langle D(V(\Delta_{k,h}u(s)) \eta^2), \Delta_{k,h}A(\cdot, s, u, Du) \rangle_{L^2(D)} ds \\ & \leq \|(T_{q,K} |\Delta_{k,h}u_0|)^{\frac{1}{2}} \eta\|_{L^2(D)}^2 \\ & \quad + (q-2^{-1}) \int_0^t \|T'_{q,K}(|\Delta_{k,h}u(s)|)^{\frac{1}{2}} |\Delta_{k,h}u(s)|^{-\frac{1}{2}} \Delta_{k,h}H(\cdot, s, Du) \eta\|_{L^2(D)}^2 ds \\ & \quad + \int_0^t \langle V(\Delta_{k,h}u(x, s)) \eta^2, \Delta_{k,h}H(\cdot, s, Du) dB_s \rangle_{L^2(D)} \end{aligned} \tag{4.11}$$

Now the second term on the left-hand side of this inequality shall be estimated. Using the decomposition introduced in (4.6) and applying Lemma 3.11, we first find for every $\varepsilon > 0$:

$$\begin{aligned} & \langle D(V(\Delta_{k,h}u(s)) \eta^2), \mathcal{A}(h) \rangle_{L^2(D)} \\ &= \kappa^{-1} \langle D(V(\Delta_{k,h}u(s)) \eta^2), D\Delta_{k,h}u \rangle_{L^2(D)} - \kappa^{-1} \langle D(V(\Delta_{k,h}u(x, s)) \eta^2), D\Delta_{k,h}u - \kappa \mathcal{A}(h) \rangle_{L^2(D)} \\ &\geq \kappa^{-1} \|D(V(\Delta_{k,h}u(s))) \cdot D\Delta_{k,h}u \eta^2\|_{L^1(D)} - 2\varepsilon \|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{-\frac{1}{2}} D\Delta_{k,h}u \eta\|_{L^2(D)}^2 \\ &\quad - c(q, \kappa, \varepsilon) \| |\Delta_{k,h}u|^q D\eta \|_{L^2(D)}^2 - \kappa^{-1} (1 - \nu^2)^{\frac{1}{2}} \|D(V(\Delta_{k,h}u(s))) \|D\Delta_{k,h}u \eta^2\|_{L^1(D)} \\ &\geq (\kappa^{-1} \mu^{\frac{1}{2}}(q) - \kappa^{-1} (1 - \nu^2)^{\frac{1}{2}} - 2\varepsilon) \|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{-\frac{1}{2}} D\Delta_{k,h}u \eta\|_{L^2(D)}^2 \\ &\quad - c(q, \kappa, \varepsilon, \|D\eta\|_{L^\infty(D)}) \|W_q(\Delta_{k,h}u)\|_{L^2(D')}^2. \end{aligned}$$

We observe from the definition of $\mu(q)$ and the second bound in (4.10) on q that the factor $\mu^{\frac{1}{2}}(q) - (1 - \nu^2)^{\frac{1}{2}}$ appearing in the previous inequality is always strictly positive. Now, for the second term in the

decomposition (4.6) we obtain via the inequalities $T''_{q,K}(t)t^2 \leq c(q)T'_{q,K}(t)t \leq c(q)T_{q,K}(t)$ on \mathbb{R}^+ and the Sobolev-Poincaré embedding (applied on every time-slice):

$$\begin{aligned}
& \left| \langle D(V(\Delta_{k,h}u(s))\eta^2), \mathcal{B}(h) \rangle_{L^2(D)} \right| \\
& \leq c(L, q) \left(\|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{-\frac{1}{2}} D\Delta_{k,h}u\eta\|_{L^2(D)} + \|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{\frac{1}{2}} D\eta\|_{L^2(D)} \right) \\
& \quad \times \|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{-\frac{1}{2}} \mathcal{B}(h)\eta\|_{L^2(D)} \\
& \leq c(L, q) \left(\|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{-\frac{1}{2}} D\Delta_{k,h}u\eta\|_{L^2(D)} + \|T_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} D\eta\|_{L^2(D)} \right) \\
& \quad \times \|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{-\frac{1}{2}} \Delta_{k,h}u\eta\|_{L^{\frac{2n}{n-2}}(D)}^{\theta} \|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{-\frac{1}{2}} \Delta_{k,h}u\eta\|_{L^2(D)}^{1-\theta} \\
& \quad \times \left(\|Du\|_{L^{\frac{2n}{(n+2)\theta}}(\text{spt } \eta)}^{\frac{2}{n+2}} + \|u\|_{L^{\frac{2}{\theta}}(\text{spt } \eta)}^{\frac{2}{n}} + \|f\|_{L^{\frac{n}{\theta}}(D)} \right) \\
& \leq \left(\|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{-\frac{1}{2}} D\Delta_{k,h}u\eta\|_{L^2(D)}^{1+\theta} + \|T_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} D\eta\|_{L^2(D)}^{1+\theta} \right) \\
& \quad \times c(n, D, T, L, q) \|T_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} \eta\|_{L^2(D)}^{1-\theta} \left(\|Du\|_{L^{\frac{2n}{(n+2)\theta}}(\text{spt } \eta)}^{\frac{2}{n+2}} + \|u\|_{L^{\frac{2}{\theta}}(\text{spt } \eta)}^{\frac{2}{n}} + \|f\|_{L^{\frac{n}{\theta}}(D)} \right)
\end{aligned}$$

for every $\theta \in (0, 1)$. We now choose $\theta = \max\{p^{-1}(\frac{n}{n+2})^2, \frac{n}{a}\}$, for which the last expression in brackets of the previous inequality is consequently bounded with probability one, according to the integrability assumptions on f, Du and the consequences on the integrability of u explained at the beginning of the proof. Young's inequality then implies

$$\begin{aligned}
\left| \langle D(V(\Delta_{k,h}u(x, s))\eta^2), \mathcal{B}(h) \rangle_{L^2(D)} \right| & \leq \|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{-\frac{1}{2}} D\Delta_{k,h}u\eta\|_{L^2(D)}^2 \\
& \quad + c(n, D, T, L, q, \|D\eta\|_{L^\infty(D)}, \varepsilon) \left(\|T_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} \eta\|_{L^2(D)}^2 + 1 \right) \\
& \quad \times \left(1 + \|Du\|_{L^{2p\frac{n+2}{n}}(\text{spt } \eta)}^{2p\frac{n+2}{n}} + \|u\|_{L^{2p(\frac{n+2}{n})^2}(\text{spt } \eta)}^{2p(\frac{n+2}{n})^2} + \|f\|_{L^a(D)}^a \right).
\end{aligned}$$

Finally, via the bounds for q in terms of n, p, a and ν , the last term in the decomposition involving $\mathcal{C}(h)$ is estimated with Young's inequality and the well-known estimates for finite difference quotients by

$$\begin{aligned}
& \left| \langle D(V(\Delta_{k,h}u(s))\eta^2), \mathcal{C}(h) \rangle_{L^2(D)} \right| \\
& \leq c \left(\|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{-\frac{1}{2}} D\Delta_{k,h}u\eta\|_{L^2(D)} + \|T_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} D\eta\|_{L^2(D)} \right) \\
& \quad \times \|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{-\frac{1}{2}} \mathcal{C}(h)\eta\|_{L^2(D)} \\
& \leq \varepsilon \|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{-\frac{1}{2}} D\Delta_{k,h}u\eta\|_{L^2(D)}^2 \\
& \quad + c(D, T, L, q, \|D\eta\|_{L^\infty(D)}, \varepsilon) \left(1 + \|Du\|_{L^{2p\frac{n+2}{n}}(\text{spt } \eta)}^{2p\frac{n+2}{n}} + \|u\|_{L^{2p(\frac{n+2}{n})^2}(\text{spt } \eta)}^{2p(\frac{n+2}{n})^2} + \|f\|_{L^a(D)}^a \right)
\end{aligned}$$

provided that $4q \leq a$. For the general case, one again has to argue more subtle, using the Sobolev embedding on time slices as for the term with $\mathcal{B}(h)$. With the analogous calculations as before this yields

$$\begin{aligned}
& \left| \langle D(V(\Delta_{k,h}u(s))\eta^2), \mathcal{C}(h) \rangle_{L^2(D)} \right| \\
& \leq \varepsilon \|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{-\frac{1}{2}} D\Delta_{k,h}u\eta\|_{L^2(D)}^2 \\
& \quad + c(n, D, T, L, q, \|D\eta\|_{L^\infty(D)}, \varepsilon) \left(\|T_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} \eta\|_{L^2(D)}^2 + 1 \right) \\
& \quad \times \left(1 + \|Du\|_{L^{2p\frac{n+2}{n}}(\text{spt } \eta)}^{2p\frac{n+2}{n}} + \|u\|_{L^{2p(\frac{n+2}{n})^2}(\text{spt } \eta)}^{2p(\frac{n+2}{n})^2} + \|f\|_{L^a(D)}^a \right).
\end{aligned}$$

It now still remains to handle the second term on the right-hand side of inequality (4.11). With the assump-

tions (2.2) on H and Young's inequality, we easily find

$$\begin{aligned} & \|T'_{q,K}(|\Delta_{k,h}u(s)|)^{\frac{1}{2}} |\Delta_{k,h}u(s)|^{-\frac{1}{2}} \Delta_{k,h}H(\cdot, s, Du) \eta\|_{L^2(D)}^2 \\ & \leq (L_H^2 + \varepsilon) \|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{-\frac{1}{2}} D\Delta_{k,h}u \eta\|_{L^2(D)}^2 \\ & \quad + c(L, L_H, \varepsilon) \left(1 + \|Du\|_{L^{2p\frac{n+2}{n}}(\text{spt } \eta)}^{2p\frac{n+2}{n}} + \|f_H\|_{L^a(D)}^a\right). \end{aligned}$$

For every $s \in (0, T)$ we now define

$$G''(u, f)(s) := \frac{2q}{c''} G_p(s) + 1 + \|Du\|_{L^{2p\frac{n+2}{n}}(D')}^{2p\frac{n+2}{n}} + \|u\|_{L^{2p(\frac{n+2}{n})^2}(D')}^{2p(\frac{n+2}{n})^2} + \|f\|_{L^a(D)}^a, \quad (4.12)$$

which is a $L^1(0, T)$ with probability one. Furthermore, we set $G_q := \frac{1}{2q} c'' G''(u, f) \geq G_p$ which immediately gives $Y_q \leq Y_p$. Then, taking into account the smallness condition (4.10), choosing ε sufficiently small and combining the previous estimates for the various terms arising in (4.11), we find a preliminary (though still K -depending) pathwise estimate

$$\begin{aligned} & \|T_{q,K}(|\Delta_{k,h}u(t)|)^{\frac{1}{2}} \eta\|_{L^2(D)}^2 + c^{-1}(L_H, \kappa, \nu) \int_0^t \|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{-\frac{1}{2}} D\Delta_{k,h}u \eta\|_{L^2(D)}^2 ds \\ & \leq \|T_{q,K}(|\Delta_{k,h}u_0|)^{\frac{1}{2}} \eta\|_{L^2(D)}^2 + c'' \int_0^t (\|T_{q,K}(|\Delta_{k,h}u(t)|)^{\frac{1}{2}} \eta\|_{L^2(D)}^2 + 1) G''(u, f) ds \\ & \quad + c \int_0^t \|f_H(s)\|_{L^a(D)}^a ds + \int_0^t \langle V(\Delta_{k,h}u(x, t)) \eta^2, \Delta_{k,h}H(\cdot, s, Du) dB_s \rangle_{L^2(D)}. \end{aligned}$$

Step 2. We may now apply in a first step Itô's formula in exactly the same way as before in the derivation of estimate (4.9):

$$\begin{aligned} & e^{-\int_0^t c'' G''(u, f) ds} \|T_{q,K}(|\Delta_{k,h}u(t)|)^{\frac{1}{2}} \eta\|_{L^2(D)}^2 \\ & + c^{-1} \int_0^t e^{-\int_0^s c'' G''(u, f) d\tilde{s}} \|T'_{q,K}(|\Delta_{k,h}u|)^{\frac{1}{2}} |\Delta_{k,h}u|^{-\frac{1}{2}} D\Delta_{k,h}u \eta\|_{L^2(D)}^2 ds \\ & \leq \|T_{q,K}(|\Delta_{k,h}u_0|)^{\frac{1}{2}} \eta\|_{L^2(D)}^2 + 1 + c \int_0^t \|f_H(s)\|_{L^a(D)}^a ds \\ & \quad + c \int_0^t e^{-\int_0^s c'' G''(u, f) d\tilde{s}} \langle V(\Delta_{k,h}u(x, t)) \eta^2, \Delta_{k,h}H(\cdot, s, Du) dB_s \rangle_{L^2(D)}. \end{aligned}$$

Step 3. Similarly to the proof of Lemma 4.2, we introduce the random time

$$\tau_R := \inf \left\{ t \in [0, T] : \int_0^t \|\Delta_{k,h}u(s)\|^{2q-1} \eta^2 |\Delta_{k,h}H(\cdot, s, Du)|\|_{L^1(D)}^2 ds > R \right\}$$

with $\tau_R = T$ when the set is empty. Differently from Lemma 4.2, the property

$$P\left(\int_0^T \|\Delta_{k,h}u(s)\|^{2q-1} \eta^2 |\Delta_{k,h}H(\cdot, s, Du)|\|_{L^1(D)}^2 ds < \infty\right) = 1$$

which is needed to have $P(\lim_{R \rightarrow \infty} \tau_R = T) = 1$ is not clear a priori. We shall prove it a posteriori.

Notice that, by Lemma 3.11,

$$\begin{aligned} & \int_0^{t \wedge \tau_R} e^{-2\int_0^s c'' G''(u, f) d\tilde{s}} \left(\int_D |V(\Delta_{k,h}u(x, s))| \eta^2 |\Delta_{k,h}H(s, Du)| dx \right)^2 ds \\ & \leq \int_0^{t \wedge \tau_R} \|\Delta_{k,h}u(s)\|^{2q-1} \eta^2 |\Delta_{k,h}H(\cdot, s, Du)|\|_{L^1(D)}^2 ds \leq R. \end{aligned}$$

Step 3a. The last calculation shows that the stochastic integral from Step 2, stopped at τ_R , is a martingale (and thus it has zero expectation). Therefore (as in Lemma 4.2)

$$\begin{aligned} & E \left[e^{-\int_0^{t \wedge \tau_R} c'' G''(u, f) ds} \left\| T_{q, K}(|\Delta_{k, h} u(t \wedge \tau_R)|)^{\frac{1}{2}} \eta \right\|_{L^2(D)}^2 \right] \\ & + c^{-1} E \left[\int_0^{t \wedge \tau_R} e^{-\int_0^s c'' G''(u, f) d\tilde{s}} \left\| T'_{q, K}(|\Delta_{k, h} u|)^{\frac{1}{2}} |\Delta_{k, h} u|^{-\frac{1}{2}} D \Delta_{k, h} u \eta \right\|_{L^2(D)}^2 ds \right] \\ & \leq \left\| T_{q, K}(|\Delta_{k, h} u_0|)^{\frac{1}{2}} \eta \right\|_{L^2(D)}^2 + 1 + c E \left[\int_0^T \|f_H(s)\|_{L^a(D)}^a ds \right]. \end{aligned}$$

At this stage we may pass to the limit $K \rightarrow \infty$ via Fatou's Lemma on the left-hand side and monotone convergence on the right-hand side, and we obtain

$$\begin{aligned} & E \left[e^{-\int_0^{t \wedge \tau_R} c'' G''(u, f) ds} \left\| W_q(\Delta_{k, h} u(t \wedge \tau_R)) \eta \right\|_{L^2(D)}^2 \right] \\ & + c^{-1} E \left[\int_0^{t \wedge \tau_R} e^{-\int_0^s c'' G''(u, f) d\tilde{s}} \left\| |\Delta_{k, h} u(s)|^{q-1} D \Delta_{k, h} u(s) \eta \right\|_{L^2(D)}^2 ds \right] \\ & \leq \left\| W_q(\Delta_{k, h} u_0) \eta \right\|_{L^2(D)}^2 + 1 + c E \left[\int_0^T \|f_H(s)\|_{L^a(D)}^a ds \right]. \end{aligned} \quad (4.13)$$

Step 3b. Next we apply Burkholder-Davis-Gundy inequality to the inequality above stopped at τ_R , raised to the power $\frac{1}{q}$. Taking the limit $K \rightarrow \infty$ as in (4.13), we get

$$\begin{aligned} & E \left[\sup_{t \in [0, T]} e^{-\frac{1}{q} \int_0^{t \wedge \tau_R} c'' G''(u, f) ds} \left\| W_q(\Delta_{k, h} u(t \wedge \tau_R)) \eta \right\|_{L^2(D)}^{\frac{2}{q}} \right] \\ & \leq \left\| W_q(\Delta_{k, h} u_0) \eta \right\|_{L^2(D)}^{\frac{2}{q}} + 1 + c E \left[\left(\int_0^T \|f_H(s)\|_{L^a(D)}^a ds \right)^{\frac{1}{q}} \right] \\ & + C E \left[\left(\int_0^{T \wedge \tau_R} e^{-2 \int_0^s c'' G''(u, f) d\tilde{s}} \left\| |\Delta_{k, h} u(s)|^{2q-1} \eta^2 |\Delta_{k, h} H(\cdot, s, Du)| \right\|_{L^1(D)}^2 ds \right)^{\frac{1}{2q}} \right]. \end{aligned}$$

Since due to Hölder's inequality we have

$$\begin{aligned} & \left\| |\Delta_{k, h} u(s)|^{2q-1} \eta^2 |\Delta_{k, h} H(\cdot, s, Du)| \right\|_{L^1(D)}^2 \\ & \leq \left\| |\Delta_{k, h} u(s)|^q \eta \right\|_{L^2(D)}^2 \left\| |\Delta_{k, h} u(s)|^{q-1} \eta |\Delta_{k, h} H(\cdot, s, Du)| \right\|_{L^2(D)}^2, \end{aligned}$$

the last term of the previous inequality, similarly to the proof of Lemma 4.2, is bounded by

$$c E [I_1^{1/2} I_2^{1/2}] \leq \frac{1}{2} E [I_1] + \frac{C^2}{2} E [I_2]$$

where

$$\begin{aligned} I_1 &= \sup_{t \in [0, T]} e^{-\frac{1}{q} \int_0^{t \wedge \tau_R} c'' G''(u, f) d\tilde{s}} \left\| W_q(\Delta_{k, h} u(t \wedge \tau_R)) \eta \right\|_{L^2(D)}^{\frac{2}{q}}, \\ I_2 &= \left(\int_0^{T \wedge \tau_R} e^{-\int_0^s c'' G''(u, f) d\tilde{s}} \left\| |\Delta_{k, h} u(s)|^{q-1} \eta |\Delta_{k, h} H(\cdot, s, Du)| \right\|_{L^2(D)}^2 ds \right)^{\frac{1}{q}}. \end{aligned}$$

Hence, we have proved that

$$\frac{1}{2} E [I_1] \leq \left\| W_q(\Delta_{k, h} u_0) \eta \right\|_{L^2(D)}^{\frac{2}{q}} + 1 + c E \left[\left(\int_0^T \|f_H(s)\|_{L^a(D)}^a ds \right)^{\frac{1}{q}} \right] + \frac{C^2}{2} E [I_2].$$

Now, by the assumptions (2.2) on H , Young's inequality and the bound on q , we have

$$\begin{aligned} \frac{C^2}{2} E[I_2] &\leq C E \left[\left(\int_0^{T \wedge \tau_R} e^{-\int_0^s c'' G''(u, f) d\tilde{s}} \left\| |\Delta_{k, h} u(s)|^{q-1} D \Delta_{k, h} u(s) \eta \right\|_{L^2(D)}^2 ds \right)^{\frac{1}{q}} \right] \\ &\quad + C + C E \left[\left(\int_0^T \|f_H(s)\|_{L^a(D)}^a ds \right)^{\frac{1}{q}} \right] + \frac{1}{4} E[I_1] \\ &\quad + C E \left[\left(\int_0^T e^{-\int_0^s c'' G''(u, f) d\tilde{s}} \|W_q(Du(s))\|_{L^2(D')}^2 ds \right)^{\frac{1}{q}} \right]. \end{aligned}$$

We observe that the last term remains bounded, due to the assumption of the lemma on the average and the choice of $G''(u, f)$ (which ensures that $c'' G''(u, f) \geq 2qG_p$). Thus, by inequality (4.13) proved above, we find

$$\frac{1}{4} E[I_1] \leq c \|W_q(\Delta_{k, h} u_0) \eta\|_{L^2(D)}^{\frac{2}{q}} + c + c E \left[\int_0^T \|f_H(s)\|_{L^a(D)}^a ds \right]^{\frac{1}{q}}$$

with a new constant.

Step 3c. In Step 3a and Step 3b we almost proved the two bounds claimed by the lemma since the previous inequality along with (4.13) gives us

$$\begin{aligned} &E \left[\sup_{t \in [0, T]} e^{-\frac{1}{q} \int_0^{t \wedge \tau_R} c'' G''(u, f) ds} \|W_q(\Delta_{k, h} u(t \wedge \tau_R)) \eta\|_{L^2(D)}^{\frac{2}{q}} \right] \\ &+ E \left[\left(\int_0^{T \wedge \tau_R} e^{-\int_0^s c'' G''(u, f) d\tilde{s}} \left\| |\Delta_{k, h} u(s)|^{q-1} D \Delta_{k, h} u(s) \eta \right\|_{L^2(D)}^2 ds \right)^{\frac{1}{q}} \right] \\ &\leq c \left(\|W_q(\Delta_{k, h} u_0) \eta\|_{L^2(D)}^2 + c + c E \left[\int_0^T \|f_H(s)\|_{L^a(D)}^a ds \right] \right)^{\frac{1}{q}}. \end{aligned}$$

It now remains to justify (as already observed above) the limit $\tau_R \rightarrow T$ as $R \rightarrow \infty$ with probability one. Indeed, since $R \mapsto \tau_R$ is non-decreasing and bounded above by T , there exists the a.s. limit

$$\tau := \lim_{R \rightarrow \infty} \tau_R$$

and $\tau(\omega) \in [0, T]$. By Fatou's lemma and monotone convergence,

$$\begin{aligned} &E \left[\left(\sup_{t \in [0, \tau]} e^{-\int_0^t c'' G''(u, f) ds} \|W_q(\Delta_{k, h} u(t)) \eta\|_{L^2(D)}^2 \right. \right. \\ &\quad \left. \left. + \int_0^\tau e^{-\int_0^s c'' G''(u, f) d\tilde{s}} \left\| |\Delta_{k, h} u(s)|^{q-1} D \Delta_{k, h} u(s) \eta \right\|_{L^2(D)}^2 ds \right)^{\frac{1}{q}} \right] \end{aligned}$$

is finite, hence the argument of the expectation is finite with probability one. Since $\int_0^T c'' G''(u, f) ds$ is finite with probability one, we get

$$\sup_{t \in [0, \tau]} \|W_q(\Delta_{k, h} u(t)) \eta\|_{L^2(D)}^2 + \int_0^\tau \left\| |\Delta_{k, h} u(s)|^{q-1} D \Delta_{k, h} u(s) \eta \right\|_{L^2(D)}^2 ds < \infty$$

with probability one. Thus (with the same inequalities used above)

$$\begin{aligned} &\int_0^\tau \left\| |\Delta_{k, h} u(s)|^{2q-1} \eta^2 |\Delta_{k, h} H(\cdot, s, Du)| \right\|_{L^1(D)}^2 ds \\ &\leq C \left(1 + \sup_{t \in [0, \tau]} \|W_q(\Delta_{k, h} u(t)) \eta\|_{L^2(D)}^2 + \int_0^\tau \left\| |\Delta_{k, h} u(s)|^{q-1} D \Delta_{k, h} u(s) \eta \right\|_{L^2(D)}^2 ds \right. \\ &\quad \left. + \int_0^T \|W_q(Du(s))\|_{L^2(D')}^2 ds + \int_0^T \|f_H(s)\|_{L^a(D)}^a ds \right)^2 < \infty \end{aligned}$$

with probability one. If $\tau(\omega) < T$, by definition of τ_R we have

$$\int_0^\tau \| |\Delta_{k,h} u(s)|^{2q-1} \eta^2 |\Delta_{k,h} H(\cdot, s, Du)| \|_{L^1(D)}^2 ds = \infty$$

which is false, hence $P(\tau = T) = 1$. Having this basic fact, the same estimates just proved give us the result of the lemma, by taking into account the inequality $|DW(\Delta_{k,h} u)| \leq q |\Delta_{k,h} u|^{q-1} |D\Delta_{k,h} u|$ and the definition of G_q (and hence of Y_q) given after (4.12). \square

5 Proof of the regularity result

Having the previous lemma at hand, we may now proceed to our main result.

Theorem 5.1. *Let u be a weak solution to the initial boundary value problem to (2.3) with initial values $u(\cdot, 0) = u_0(\cdot) \in W^{1,a-2}(D, \mathbb{R}^N)$. Assume further the assumptions (2.1) with $\nu > (n-2)/n$ such that*

$$L_H^2 < (L_H^*)^2(n) := \frac{2}{\kappa(n-1)} \left(\left[1 - \left(\frac{n-2}{n} \right)^2 \right]^{\frac{1}{2}} - [1 - \nu^2]^{\frac{1}{2}} \right).$$

Then there exists $\alpha > 0$ depending only on n, ν and a such that for every subset $D_c \Subset D$ we have

$$P(\|u\|_{C^{0,\alpha}(D_c \times [0,T], \mathbb{R}^N)} < \infty) = 1.$$

Proof. To prove the result, we want to apply Proposition 3.9. Therefore, the crucial point is to show higher integrability of Du for “great” powers with probability one, in order that hypothesis (3.6) of the proposition is satisfied. We start by defining a sequence

$$\begin{aligned} \tilde{q}_0 &:= 1, \\ \tilde{q}_{j+1} &:= \min \left\{ q_j \frac{n+2}{n}, 1 + q_j \frac{n+2}{n} \frac{a-4}{a}, \frac{a-2}{2}, q_j + 1 \right\} \text{ for } j \geq 1. \end{aligned}$$

Before defining a further sequence (q_j) in order to perform the iteration, we make some observations on $L_H^*(s)$ as a function in $s \in [1, 2/(1-\nu)]$ (we note that $L_H^*(2q)$ already appeared in hypothesis (4.10) which gave an upper bound for q in the iteration). Clearly, $L_H^*(s)$ is strictly decreasing in s , with $L_H^*(2/(1-\nu)) = 0$.

We now set $q_j = \tilde{q}_j$ as long as $\tilde{q}_j > \tilde{q}_{j-1}$ and $L_H^*(2\tilde{q}_j) > L_H$, and for the first index j which doesn't satisfy these assumptions any more we set $q_j = q^*$ for a number $q^* > n/2$ (which is determined below). In what follows we shall denote this set of indices by $J \subset N_0$. We first study some properties of the sequence \tilde{q} and give a definition of the final member q^* of the sequence $(q_j)_{j \in J}$: the first and the forth term in the rewritten formula for \tilde{q} are strictly increasing in j and diverge for $j \rightarrow \infty$, whereas the monotonicity properties of the second term depend on both the values of a and the size of q_j . More precisely, if $a \geq 2(n+2)$, then the second term increases with j and diverges for $j \rightarrow \infty$, but for every $a \in (n+2, 2(n+2))$ it increases only up to $q_{max}(a, n) = na/(4(n+2) - 2a) > n/2$. Observing $L_H^*(n) > L_H$ by assumption, we thus define

$$q^* := \text{arbitrary number in } \left(\frac{n}{2}, \min \left\{ (L_H^*)^{-1}(L_H), \frac{a-2}{2}, q_{max} \right\} \right).$$

It is easy to calculate that this number q^* is reached after a finite number of steps (depending only on n, ν, a and the difference $q_{max} - q^*$ (in the sense that the number of steps diverges as $q^* \nearrow q_{max}$), hence $|J| < \infty$, i.e. $(q_j)_{j \in J}$ is a finite sequence.

We are now going to establish by induction that for every $j \in J$ we have

- (i) $\sup_{|h| < \text{dist}(D_j, \partial D_{j-1})} E[\|Y_{q_j}^{q_j} W_{q_j}(\Delta_{k,h} u)\|_{V^2(D_j \times (0,T))}^{2/q_j}] \leq C_j$ for all $k \in \{1, \dots, n\}$,
- (ii) $\sup_{|h| < \text{dist}(D_j, \partial D_{j-1})} E[\|Y_{q_j}^{q_j} W_{q_j}(\Delta_{k,h} u)\|_{L^2 \frac{2+2}{n}(D_j \times (0,T))}^{2/q_j}] \leq c(n, D_j) C_j$ for all $k \in \{1, \dots, n\}$,
- (iii) $E[\|Y_{q_j}^{q_j} W_{q_j}(Du)\|_{L^2 \frac{n+2}{n}(D_j \times (0,T))}^{2/q_j}] \leq \tilde{C}_j$,
- (iv) $Du \in L^\infty(0, T; L^{2q_j}(D_j, \mathbb{R}^{nN}))$ with probability one.

Here $(Y_{q_j})_{j \in J}$ is a sequence of random variables given by $Y_{q_j}(t, \omega) = \exp(-\int_0^t G_{q_j}(s, \omega) ds)$ for each $j \in J$, for a sequence of functions $(G_{q_j})_{j \in J}$ which are in $L^1(0, T)$ with probability one and which will be determined later, and $(D_j)_{j \in J}$ is a monotone decreasing sequence of open sets satisfying $D_c \subset D_j \subset D_{j-1} \subset \dots \subset D_0 \subset D_{-1} = D$.

We start by setting

$$Y_1 := e^{-\frac{1}{2} \int_0^t c' G'(u, f) ds},$$

where $G'(u, f)$ was defined in (4.7). It is obvious from its definition that $Y_1: [0, T] \times \Omega \rightarrow (0, 1]$ satisfies $P(\inf_{t \in [0, T]} Y_1 > 0) = 1$. We then observe from Lemma 4.2 that

$$\begin{aligned} \sup_{|h| < d} E \left[\sup_{t \in (0, T)} \|Y_1 \Delta_{k, h} u\|_{L^2(D_0)}^2 + \int_0^T \|Y_1 D \Delta_{k, h} u\|_{L^2(D_0)}^2 dt \right] \\ \leq c' \left(\|D_k u_0\|_{L^2(D)}^2 + 1 + E[\|f_H^{\frac{a}{a-2}}\|_{L^2(D_T)}^2] \right) =: C_0 \end{aligned}$$

is satisfied for every open set D_0 compactly supported in D . By definition of the space V^2 , this establishes the statement (i)₀. Furthermore, (ii)₀ follows immediately from the Sobolev embedding (2.4), applied for P -almost every ω to the functions $Y_1 \Delta_{k, h} u$, for $k \in \{1, \dots, n\}$. To conclude the first step of the iteration it only remains to justify the statements (iii)₀ and (iv)₀. To this end we take advantage of Theorem 3.6 twice, in the way as explained in Remark 3.7 (and actually as already performed in Corollary 4.4). First we apply it with the choices $p = q = 2q_0 \frac{n+2}{n}$ to the inequality from (ii)₀ (for all $k \in \{1, \dots, n\}$), leading to the existence of Du in the Lebesgue space $L^{2(n+2)/n}(D_0 \times (0, T), \mathbb{R}^{nN})$ with the required estimate for the average of $Y_1 Du$; secondly, we apply it with the choice $p = 2q_0$ and $q = \infty$ to (i)₀ – more precisely to the first term in the V^2 -norm – and, keeping in mind the pathwise strict positivity of Y_1 , we end up with the existence of Du in $L^\infty(0, T; L^2(D_0, \mathbb{R}^{nN}))$ with probability one.

We now proceed to the inductive step. Assume for a given $j \in J$ that (i)_ℓ–(iv)_ℓ are valid on open sets $D_\ell \subset D_{\ell-1}$ with random variables $Y_{q_\ell}: [0, T] \times \Omega \rightarrow (0, 1]$ of the required form for all $\ell \in \{0, \dots, j-1\}$. Then, keeping in mind (iii)_{j-1} and the definition of the number q^* , we note that the assumptions of Lemma 4.5 are satisfied (for p, D' replaced by q_{j-1}, D_{j-1}), and we hence deduce (with the admissible choice $q = q_j$) the estimate

$$\begin{aligned} \sup_{|h| < d_j} E \left[\left(\sup_{t \in (0, T)} \|Y_{q_j}^{q_j} W_{q_j}(\Delta_{k, h} u)\|_{L^2(D_j)}^2 + \int_0^T \|Y_{q_j}^{q_j} D W_{q_j}(\Delta_{k, h} u)\|_{L^2(D_j)}^2 dt \right)^{\frac{1}{q_j}} \right] \\ \leq c \left(\|W_{q_{j-1}}(D_k u(x, 0))\|_{L^2(D)}^2 + 1 + E[\|f_H(s)\|_{L^a(D_T)}^a] \right)^{\frac{1}{q_j}} =: C_j \end{aligned}$$

for every $k \in \{1, \dots, n\}$, a domain $D_j \subset D_{j-1}$ satisfying $d_j := \text{dist}(D_j, \partial D_{j-1}) > 0$ and a random variable Y_{q_j} defined via G_{q_j} given in Lemma 4.5 and satisfying in particular $P(\inf_{t \in [0, T]} Y_{q_j} > 0) = 1$. This shows (i)_j, and (ii)_j in turn is an immediate consequence after the application of the Sobolev embedding as above. Moreover, the statements (iii)_j and (iv)_j again follow from (ii)_j and (i)_j, respectively, after the application of Theorem 3.6 with the choices $p = q = 2q_j \frac{n+2}{n}$ and $p = 2q_j, q = \infty$, respectively. This finishes the proof of the induction.

As an immediate consequence of the induction, we can now conclude the desired higher integrability result to a great power, via the following observation. Via (iv) we find in the limit

$$Du \in L^\infty(0, T; L^{2q^*}(D_c, \mathbb{R}^{nN}))$$

with probability one, and by definition the exponent $2q^*$ is greater than the space dimension n . Hence, assumption (3.6) of Proposition 3.9 is guaranteed. For its application we still need to check the integrability condition on $a(x, s), b(x, s)$ given by

$$a(x, s) := \text{div } A(x, s, u, Du) \quad \text{and} \quad b(x, s) := H(x, s, Du).$$

Since $A(x, t, u, z)$ is differentiable in x, u , and z with bounds (2.1), we obtain $a \in L^2(D_c \times (0, T), \mathbb{R}^N)$ with probability one as a direct consequence of $Du \in V^2(D_c \times (0, T), \mathbb{R}^{nN})$ and $f \in L^a(D_T) \subset L^4(D_T)$.

Furthermore, the growth of H according to (2.2) with $f_H \in L^a(D_T \times \Omega)$ implies $b \in L^{2+\varepsilon}(0, T; L^2(D_c, \mathbb{R}^{n'N}))$ with probability one. Thus, Proposition 3.9 yields the asserted Hölder continuity of u with probability one and finishes the proof of the theorem. \square

6 Regularity of the average due to noise

It has been recently proved that a Stratonovich bilinear multiplicative noise may have a regularizing effect on certain classes of PDEs, see [8] for a review, based on a number of works including [9, 11, 1]. In most cases, *uniqueness* by noise is the topic of these works. The problem of the interaction between noise and *singularities* is more difficult and less explored. But two examples are known:

- (i) for linear transport equations of the form

$$du = (b(x, t) \cdot Du) dt + \sigma Du \circ dB_t$$

with $b \in C(0, T; C_b^\alpha(\mathbb{R}^n, \mathbb{R}^n))$, where regular initial condition may develop discontinuities in finite time in the case $\sigma = 0$ (think of the simple example in dimension $n = 1$ given by $b(x) = -\text{sign}(x)\sqrt{|x|}$), it is known that C^1 -smoothness is preserved for $\sigma \neq 0$, see [9, 10], where similar results have been also proved for linear continuity equations;

- (ii) for the point vortex motion associated to the 2D Euler equations, it has been proved that coalescence of vortices cannot happen when a suitable Stratonovich bilinear multiplicative noise is added to the equations, see [11].

One should also notice that other singularities, like those arising in the inviscid Burgers equation, do not disappear under noise, see [8], so each equation requires its own understanding and investigation. Moreover, no general method exists to investigate these kind of properties.

Our aim here is to give a simple partial result in this direction (namely the effect of noise on singularities) for *linear systems*. We consider the linear stochastic system with Stratonovich bilinear multiplicative noise of the form

$$du = \text{div} (A(x, t) Du) dt + \sigma Du \circ dB_t, \quad u|_{t=0} = u_0 \quad (6.1)$$

with bounded measurable coefficient matrix A , where B_t is a Brownian motion in \mathbb{R}^n , defined on a filtered probability space $(\Omega, \mathcal{F}_t, P)$. The space variable x varies in a possibly unbounded regular open domain $D \subset \mathbb{R}^n$. On A we assume that there exist $\lambda_0, \lambda_1 > 0$ such that

$$\lambda_0 |\xi|^2 \leq \langle A(x, t) \xi, \xi \rangle \quad \text{and} \quad |A(x, t) \xi| \leq \lambda_1 |\xi| \quad (6.2)$$

for all $\xi \in \mathbb{R}^{nN}$, a. e. $(x, t) \in D \times [0, T]$. Actually, this is analogous to assumption (2.1)₂ (then ν corresponds to the ratio $\frac{\lambda_0}{\lambda_1}$), rewritten for vector fields which are linear in the gradient variable. We further note that for now we do not assume any regularity with respect to x , but at the same time we do not allow any dependency on Ω . Let us clarify the vector notation used in the stochastic part: $\sigma Du(x, t) \circ dB_t$ is a vector with N components, and

$$(\sigma Du(x, t) \circ dB_t)^\alpha = \sigma \sum_{i=1}^n D_i u^\alpha(x, t) \circ dB_t^i.$$

Remark 6.1. *Let us recall that Stratonovich noise is the natural one for modelling: the so called Wong-Zakai principle, proved for several classes of SPDEs (see for instance the appendix of [9] for the linear transport equation), states that solutions $u_n(x, t)$ of deterministic equations with smooth random coefficients $B_n(t)$ of the form*

$$\frac{\partial u_n}{\partial t} = \text{div} (A(x, t) Du_n) + \sigma Du_n \frac{dB_n(t)}{dt}, \quad u_n|_{t=0} = u_0$$

converge (in proper topologies and under proper assumptions on B_n , the details depend on the problem and result) to solutions u of the previous SPDE with Stratonovich noise (not Itô noise). We have stated the

principle for our system of parabolic equations just for sake of definiteness, but in fact it has not been proved before in this generality. We do not want to give a proof here, which would require a considerable work. We only quote this fact by analogy with other equations, as a general motivation for the choice of Stratonovich noise.

Let us give the definition of weak solution to equation (6.1), similarly to [9]. To understand one of the requirements (the fact that $s \mapsto \int_D u(x, s) D\varphi(x) dx$ must have a modification which is a continuous adapted semi-martingale), we recall a few facts about Stratonovich stochastic integrals, taken for instance from [17]. If B_t is a (Ω, F_t, P) -Brownian motion in \mathbb{R}^n and $X(t)$ is a continuous F_t -adapted semi-martingale, the following uniform-in-time limit exists in probability

$$\int_0^t X(s) \circ dB_s = \lim_{n \rightarrow \infty} \sum_{t_i \in \pi_n, t_i \leq t} \frac{X(t_{i+1} \wedge t) + X(t_i)}{2} (B_{t_{i+1} \wedge t} - B_{t_i})$$

and is called Stratonovich integral of X with respect to B . Here π_n is a sequence of finite partitions of $[0, T]$ with size $|\pi_n| \rightarrow 0$ and elements $0 = t_0 < t_1 < \dots$. Under the same assumptions it is defined the joint quadratic variation between X and B :

$$[X, B]_t = \lim_{n \rightarrow \infty} \sum_{t_i \in \pi_n, t_i \leq t} (X(t_{i+1} \wedge t) - X(t_i)) (B_{t_{i+1} \wedge t} - B_{t_i}),$$

and they are related to the Itô integral

$$\int_0^t X(s) dB_s = \lim_{n \rightarrow \infty} \sum_{t_i \in \pi_n, t_i \leq t} X(t_i) (B_{t_{i+1} \wedge t} - B_{t_i})$$

(which is defined under more general assumptions on X) by the formula

$$\int_0^t X(s) \circ dB_s = \int_0^t X(s) dB_s + \frac{1}{2} [X, B]_t.$$

Definition 6.2. If $u_0 \in L^2_{loc}(D, \mathbb{R}^N)$, we say that a random field $u(x, t)$ is a weak solution of equation (6.1) if:

- (i) with probability one, we have $u \in V^2(B_T, \mathbb{R}^N)$ for all bounded open sets $B \subset D$, where $B_T = B \times (0, T)$,
- (ii) for all $\varphi \in C_0^\infty(D, \mathbb{R}^N)$, the \mathbb{R}^n -valued process $s \mapsto \int_D u(x, s) D\varphi(x) dx$ has a modification which is a continuous adapted semi-martingale, and for all $t \in [0, T]$, we have P -a. s.

$$\begin{aligned} \int_D u(x, t) \varphi(x) dx + \int_0^t \int_D A(x, s) Du(x, s) D\varphi(x) dx ds + \sigma \int_0^t \left(\int_D u(x, s) D\varphi(x) dx \right) \circ dB_s \\ = \int_D u_0(x) \varphi(x) dx. \end{aligned}$$

A posteriori, from the equation itself, it follows that for all $\varphi \in C_0^\infty(D, \mathbb{R}^N)$ the real-valued process $s \mapsto \int_D u(x, s) \varphi(x) dx$ has a continuous modification. We shall always use it. Notice further that we give the following meaning to the vector notation above:

$$\int_0^t \left(\int_D u(x, s) D\varphi(x) dx \right) \circ dB_s = \sum_{i=1}^n \sum_{\alpha=1}^N \int_0^t \left(\int_D u^\alpha(x, s) D_i \varphi^\alpha(x) dx \right) \circ dB_s^i.$$

Proposition 6.3. A weak solution in the previous Stratonovich sense satisfies the Itô equation

$$\begin{aligned} \int_D u(x, t) \varphi(x) dx + \int_0^t \int_D A(x, s) Du(x, s) D\varphi(x) dx ds + \sigma \int_0^t \left(\int_D u(x, s) D\varphi(x) dx \right) dB_s \\ = \int_D u_0(x) \varphi(x) dx + \frac{\sigma^2}{2} \int_0^t \int_D u(x, s) \Delta \varphi(x) dx ds \end{aligned}$$

for all $\varphi \in C_0^\infty(D, \mathbb{R}^N)$. The converse is also true. With a language similar to that of Definition 6.2, we could say that u is a weak solution of the Itô equation

$$du = \operatorname{div} \left(\left(A(x, t) + \frac{\sigma^2}{2} \right) Du \right) dt + \sigma Du dB_t, \quad u|_{t=0} = u_0. \quad (6.3)$$

Proof. From the facts recalled above about Stratonovich integrals we have

$$\begin{aligned} \int_0^t \left(\int_D u^\alpha(x, s) D_i \varphi^\alpha(x) dx \right) \circ dB_s^i &= \int_0^t \left(\int_D u^\alpha(x, s) D_i \varphi^\alpha(x) dx \right) dB_s^i \\ &\quad + \frac{1}{2} \left[\int_D u^\alpha(x, \cdot) D_i \varphi^\alpha(x) dx, B^i \right]_t. \end{aligned}$$

Hence, we get

$$\begin{aligned} \int_D u(x, t) \varphi(x) dx + \int_0^t \int_D A(x, s) Du(x, s) D\varphi(x) dx ds + \sigma \int_0^t \left(\int_D u(x, s) D\varphi(x) dx \right) dB_s \\ = \int_D u_0(x) \varphi(x) dx - \frac{\sigma}{2} \sum_{i=1}^n \sum_{\alpha=1}^N \left[\int_D u^\alpha(\cdot) D_i \varphi^\alpha(x) dx, B^i \right]_t. \end{aligned}$$

By the equation in Definition 6.2 we also have

$$\begin{aligned} \int_D u(x, t) D_i \varphi(x) dx + \int_0^t \int_D A(x, s) Du(x, s) DD_i \varphi(x) dx ds \\ = \int_D u_0(x) D_i \varphi(x) dx - \sigma \int_0^t \left(\int_D u(x, s) DD_i \varphi(x) dx \right) \circ dB_s. \end{aligned}$$

Moreover, recall that

$$\int_0^t \left(\int_D u(x, s) DD_i \varphi(x) dx \right) \circ dB_s = \sum_{j=1}^n \int_0^t \left(\int_D u(x, s) D_j D_i \varphi(x) dx \right) \circ dB_s^j.$$

Thus, by the classical rules about quadratic variation, see [17], we have

$$\begin{aligned} \sum_{\alpha=1}^N \left[\int_D u^\alpha(x, \cdot) D_i \varphi^\alpha(x) dx, B^i \right]_t &= \left[\int_D u(x, \cdot) D_i \varphi(x) dx, B^i \right]_t \\ &= -\sigma \int_0^t \left(\int_D u(x, s) D_i D_i \varphi(x) dx \right) ds. \end{aligned}$$

The proof that the Stratonovich equation yields the Itô one is complete, and the proof of the converse statement is the same (recall that the existence of the continuous modification in (ii) of Definition 6.2 follows immediately from the equation in Definition 2.2). \square

The degree of parabolicity of the Itô SPDE (6.3) is the same as the one of (6.1), it is given just by the properties of $A(x, t)$. The term $\frac{\sigma^2}{2} \Delta u(x, t) dt$ is fully compensated by the Itô term $\sigma Du(x, t) dB_t$ and does not contribute to any additional parabolicity. This is a well recognized phenomenon in the theory of SPDEs, see for instance [16]. A simple way to see this fact is to consider the case $A \equiv 0$.

Proposition 6.4. *Consider the equation*

$$du = Du \circ dB_t, \quad u|_{t=0} = u_0 \quad (6.4)$$

in the full space $D = \mathbb{R}^n$, where B is an n -dimensional Brownian motion and $u : D \times [0, T] \times \Omega \rightarrow \mathbb{R}^N$. This is equivalent (when formulated in a weak sense) to the equation

$$du = \frac{1}{2} \Delta u dt + Du dB_t, \quad u|_{t=0} = u_0. \quad (6.5)$$

Assume $u_0 \in L^2(D, \mathbb{R}^N)$. Then

$$u(x, t) = u_0(x + B_t)$$

is a weak solution, in the sense that

(i') with probability one, we have $\int_0^T \int_{\mathbb{R}^n} |u(x, t)|^2 dx dt < \infty$,

(ii') condition (ii) of Definition 6.2 hold true.

Proof. Condition (i') comes from

$$\int_0^T \int_{\mathbb{R}^n} |u(x, t)|^2 dx dt = \int_0^T \int_{\mathbb{R}^n} |u_0(x + B_t)|^2 dx dt = \int_0^T \int_{\mathbb{R}^n} |u_0(x)|^2 dx dt < \infty.$$

Condition (ii) of Definition 6.2 is due to the following argument. For every $\psi \in C_0^\infty(\mathbb{R}^n, \mathbb{R}^N)$ we have

$$\int_{\mathbb{R}^n} u(x, t) \psi(x) dx = \int_{\mathbb{R}^n} u_0(x + B_t) \psi(x) dx = \int_{\mathbb{R}^n} u_0(x) \psi(x - B_t) dx$$

and $\psi(x - B_t)$ is the semi-martingale

$$\psi(x - B_t) = \psi(x) - \int_0^t D\psi(x - B_s) dB_s + \frac{1}{2} \int_0^t \Delta\psi(x - B_s) ds.$$

Consequently, we obtain

$$\begin{aligned} \int_{\mathbb{R}^n} u(x, t) \psi(x) dx &= \int_{\mathbb{R}^n} u_0(x) \psi(x) dx - \int_0^t \left(\int_{\mathbb{R}^n} u_0(x) D\psi(x - B_s) dx \right) dB_s \\ &\quad + \frac{1}{2} \int_0^t \left(\int_{\mathbb{R}^n} u_0(x) \Delta\psi(x - B_s) dx \right) ds, \end{aligned}$$

which shows that the stochastic process $s \mapsto \int_{\mathbb{R}^n} u(x, s) \psi(x) dx$ has a modification which is a continuous adapted semi-martingale. In addition, this computation may also be used to prove the equivalence with the Itô formulation (6.5). Finally,

$$\begin{aligned} \int_{\mathbb{R}^n} u(x, t) \varphi(x) dx - \int_{\mathbb{R}^n} u_0(x) \varphi(x) dx + \int_0^t \left(\int_{\mathbb{R}^n} u(x, s) D\varphi(x) dx \right) \circ dB_s \\ = \int_{\mathbb{R}^n} u_0(x) \varphi(x - B_t) dx - \int_{\mathbb{R}^n} u_0(x) \varphi(x) dx + \int_0^t \left(\int_{\mathbb{R}^n} u_0(x) D\varphi(x - B_s) dx \right) \circ dB_s, \end{aligned}$$

and this is equal to zero because

$$\varphi(x - B_t) = \varphi(x) - \int_0^t D\varphi(x - B_s) \circ dB_s.$$

Thus, also condition (iii) is satisfied, and the proof is complete. \square

Remark 6.5. The previous proposition shows that the Itô equation (6.5) has no regularizing properties, in spite of the presence of the term $\frac{1}{2} \Delta u$ (it is fully compensated by the Itô term). In particular, if $u_0 = 1_{x_1 > 0}$, the solution $u(x, t) = 1_{B_t^1 < x_1}$ is discontinuous in x for every given (t, ω) . At the same time

$$E[u(x, t)] = E[u_0(x + B_t)] = P(B_t^1 < x_1)$$

is smooth. This means, there are easy examples of a weak solution which have a smooth average, but which are irregular with probability one. Thus, smoothness of $E[u(x, t)]$ does not imply smoothness of $u(x, t)$, and so in general the regularity of $E[u(x, t)]$ is not enough to hope for regularity of $u(x, t)$ itself. However, it is important to observe that this example started from an irregular initial data, and that this singularity was preserved in time. Obviously, the same reasoning applies to see that in fact every solution to (6.4) is Hölder continuous in D_T if u_0 is additionally assumed to be Hölder continuous, so in particular in the case $u_0 \in W^{1,q}(D, \mathbb{R}^N)$ for some $q > n$ (which was always required for the regularity statements before).

Now let us come back to weak solutions to the general linear system (6.1) with Stratonovich noise. The crucial observation is that the average of u solves an equation with improved parabolicity. Let us first recall the classical definition used also before in this paper. If $v_0 \in L^2_{loc}(D, \mathbb{R}^N)$, we say that a (deterministic) function $v(x, t)$ is a weak solution of the parabolic equation

$$\frac{\partial v}{\partial t} = \operatorname{div} \left(\left(A(x, t) + \frac{\sigma^2}{2} \right) Dv \right), \quad v|_{t=0} = v_0 \quad (6.6)$$

if $v \in V^2_{loc}(D_T, \mathbb{R}^N)$ (in the sense of (i) of Definition 6.2 above) and

$$\begin{aligned} \int_D v(x, t) \varphi(x) dx + \int_0^t \int_D A(x, s) Dv(x, s) D\varphi(x) dx ds \\ = \int_D v_0(x) \varphi(x) dx - \frac{\sigma^2}{2} \int_0^t \int_D Dv(x, s) D\varphi(x) dx ds \end{aligned}$$

for all $\varphi \in C_0^\infty(D, \mathbb{R}^N)$.

Proposition 6.6. *If u is a weak solution of equation (6.1), then*

$$v(x, t) := E[u(x, t)]$$

is a weak solution of the parabolic equation (6.6).

Proof. Step 1. We first observe that for these linear systems, we have the a priori boundedness of the solution u in the sense that

$$\sup_{t \in [0, T]} E \left[\int_B |u(x, t)|^2 dx \right] + E \left[\int_0^T \int_B |Du(x, t)|^2 dx dt \right] < \infty \quad (6.7)$$

for all bounded sets $B \subset D$. The proof of this property follows the line of arguments of the proof of Lemma 4.1 (but is in fact much easier). We do not want to go into details, but only mention the peculiarities. First, by the linear structure in (6.1) the function G_0 appearing in Lemma 4.1 can be chosen constant. This explains, why the estimate (6.7) doesn't involve weights as before. Furthermore, in Lemma 4.1 we were content with a bound for the expected value of the spatial derivatives of u only. However, adjusting the arguments from Step 3b in the proof of Lemma 4.2, we obtain a bound for the average for the full V^2 -norm for every bounded set B compactly supported in D . This immediately gives (6.7).

Step 2. The regularity property $v \in L^\infty(0, T; L^2_{loc}(D, \mathbb{R}^N))$ is a direct consequence of the first condition in (6.7) from Step 1. In order to prove that also $v \in L^2(0, T; W^{1,2}_{loc}(D, \mathbb{R}^N))$ holds true, we first observe that we have (a. s. in t)

$$E \left[\int_D D_i u^\alpha(x, s) \psi(x) dx \right] = -E \left[\int_D u^\alpha(x, s) D_i \psi(x) dx \right]$$

for all $\psi \in C_0^\infty(D, \mathbb{R})$. This implies (by the integrability derived in (6.7))

$$\int_D E[D_i u^\alpha(x, s)] \psi(x) dx = - \int_D E[u^\alpha(x, s)] D_i \psi(x) dx,$$

which in turn gives us that $E[u^\alpha(x, s)]$ is weakly differentiable in x with partial derivative equal to $E[D_i u^\alpha(x, s)]$. Thus, v is weakly differentiable in x and

$$\begin{aligned} \int_0^T \int_B |Dv(x, t)|^2 dx dt &= \int_0^T \int_B |DE[u(x, t)]|^2 dx dt = \int_0^T \int_B |E[Du(x, t)]|^2 dx dt \\ &\leq \int_0^T \int_B E[|Du(x, t)|^2] dx dt = E \left[\int_0^T \int_B |Du(x, t)|^2 dx dt \right] < \infty \end{aligned}$$

for all bounded $B \subset D$. The regularity properties of v have been checked.

Step 3. The property $E[\int_0^T \int_B |u(x, t)|^2 dx dt] < \infty$ implies that the Itô integral in the equation of Proposition 6.3 is a martingale, hence it has zero expected value. By the same assumption, we can interchange expectation and integrals, and we get

$$\begin{aligned} \int_D v(x, t) \varphi(x) dx + E \left[\int_0^t \int_D A(x, s) Du(x, s) D\varphi(x) dx ds \right] \\ = \int_D u_0(x) \varphi(x) dx + \frac{\sigma^2}{2} \int_0^t \int_D v(x, s) \Delta \varphi(x) dx ds. \end{aligned}$$

From the property $E[\int_0^T \int_B |Du(x, t)| dx dt] < \infty$ and the boundedness of A it follows that

$$E \left[\int_0^t \int_D A(x, s) Du(x, s) D\varphi(x) dx ds \right] = \int_0^t \int_D A(x, s) E[Du(x, s)] D\varphi(x) dx ds.$$

Since we know from Step 2 that $E[Du(x, s)] = Dv(x, s)$, the proof is complete. \square

Let us now explain the possibly regularizing effect of noise. Assume $\sigma = 0$. Then, as already explained in the introduction, weak solutions may miss full regularity. One can find in [26] an example of matrix A satisfying assumption (6.2) and an example of a weak solution to the associated parabolic system (1.2) such that it is Hölder continuous on a local time interval and then its L^∞ norm blows-up. More precisely, this matrix turns out to have an ellipticity ratio $\frac{\lambda_1}{\lambda_0}$ which is smaller than the critical one employed before, which was an essential ingredient in order to obtain globally Hölder continuous weak solutions (see [14, 13, 15]). However, the matrix constructed by Stará and John [26] also fails to satisfy the regularity with respect to x , i.e. the matrix A is not differentiable in x . For this reason it is not clear whether the counterexample could be constructed due to the small ellipticity ratio or the low regularity in x or a combination of both. As far as we know there is no counterexample available in the literature which answers this question, and so even in the deterministic setting this irregularity phenomenon for weak solutions of parabolic systems is not understood completely. Instead, for the elliptic (stationary) case Koshelev was able to give a sharp result, namely that (in the linear case considered in this section) full Hölder continuity of the weak solution to $\operatorname{div}(A(x)Du) = 0$ holds provided that the matrix A is symmetric (for simplicity), measurable, bounded, and satisfies (6.2) with

$$\frac{\lambda_1 - \lambda_0}{\lambda_1 + \lambda_0} \sqrt{1 + \frac{(n-2)^2}{n-1}} < 1.$$

The sharpness of this condition follows by a modification of De Giorgi's famous counterexample [5], see [15, Section 2.5]. Returning to the parabolic setting we now state a consequence from the previous Proposition 6.6, which for randomly perturbed systems (6.1) gives a regularity result for the average $E[u(x, t)]$ if the matrix is assumed to be regular with respect to x . Since the existence of a deterministic counterexample is not clear, the Stratonovich multiplicative noise is only possibly regularizing, but in any case it might be of its own interest since the noise improves the parabolicity of the equation solved by the average.

Proposition 6.7. *Assume $q > n$, A with property (6.2) such that $|D_x A|$ is bounded uniformly by some constant $L > 0$, and let $D \subset \mathbb{R}^n$ be a bounded, regular domain. Then there exists $\sigma_0 \geq 0$ such that for all $\sigma > \sigma_0$, all initial conditions $u_0 \in W^{1,q}(D, \mathbb{R}^N)$, and all weak solutions u of equation (6.1) satisfying (6.7), we have that $(x, t) \mapsto E[u(x, t)]$ is locally Hölder continuous on $D \times [0, T]$. One can take σ_0 depending only on $\frac{\lambda_1}{\lambda_0}$.*

Proof. The matrix $A(x, t) + \frac{\sigma^2}{2}I$ satisfies the assumptions needed for the deterministic regularity results in [14, 13, 15], see also Theorem 1.1, for all σ greater than some σ_0 which can be defined in terms of $\frac{\lambda_1}{\lambda_0}$. This implies that any weak solution v of equation (6.6) is locally Hölder continuous on $D \times [0, T]$. It is sufficient to apply this result to $v(x, t) = E[u(x, t)]$. \square

Remark 6.8. Given the ratio $\frac{\lambda_1}{\lambda_0}$, the result is true for all matrices A with that ratio and all (regular) initial conditions. Intuitively speaking it looks impossible that regularization comes from the operation of mathematical expectation: it could regularize problems with special symmetries, such that singularities for different ω 's average out (compare Remark 6.5 for this phenomenon under an irregular initial condition). But here A and u_0 are quite generic (though regular). Thus we believe that Hölder regularization takes place at the level of u itself. However, this problem is open.

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